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## Strong almost convergence of order (r) of the Fourier series

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1. Let $f(t)$ be a periodic function with period $2 \pi$ and be integrable in the sense of Lebesgue over $(-\pi, \pi)$.
Let the Fourier series associated with the function $f(t)$ be

$$
f(t)=\frac{1}{2} a_{0}+\sum_{n=1}^{\infty}\left(a_{n} \cos n t+b_{n} \sin n t\right) .
$$

We write
and

$$
\phi(t)=\frac{1}{2}\{f(x+t)+f(x-t)-2 S\}
$$

$$
F(t)=\int_{0}^{t} f(u) d u
$$

2. In 1948, Lorentz [2] discussed a new process of summation which assigns a general limit Lim $x_{n}$ to certain bounded sequences $x=\left\{x_{n}\right\}$. This method is analogous to the mean values which are used in the theory of almosc periodic functions, furthermore, it is narrowly connected with limit of Banach [1]. The sequences which are summable by this method are called almost convergent sequences.

In this paper we shall first define a new method of summation which will be called "strong almost convergent" and also we shall find the sufficient conditions which sum the Fourier series by this method.
Definition: An infinite series $\Sigma a_{n}$ with the sequence of partial sum $S_{n}$ is said to be strongly almost convergent of any positive order $r$ to $s$., if

$$
\sum_{v=0}^{n}\left|S_{v+p}-S\right|^{r}=0(n), \text { as } n \rightarrow \infty
$$

uniformly with respect to $p$.
The object of this paper is to prove the following.
Theorem : If $p$ and $q$ are conjugate indices, $p>1$.

$$
\int_{0}^{t}|\phi(u)|^{p} d u=O\left(t \cdot \log ^{p / q} 1 / t\right)
$$

and

$$
\int_{0}^{t} \phi(u) d u=o\left(t / \log \frac{1}{t}\right) \text {, as } t \rightarrow 0
$$

Then

$$
\sum_{\nu=0}^{n}\left|S_{\nu+k}-S\right| r=o(n), \text { as } n \rightarrow \infty,
$$

Uniformly with respect to $k$.
3. Proof of the Theorem : With the usual standard simplifications, we write,

$$
\begin{aligned}
& S_{n+k}=\frac{2}{\pi} \int_{0}^{\pi} f(u) \frac{\operatorname{Sin}\left(n+k+\frac{1}{2}\right) u}{2 \operatorname{Sin} \frac{1}{2} u} d u \\
& =\frac{2}{\pi}\left\{\int_{0}^{\eta_{1}}+\int_{\eta_{1}}^{\eta_{2}}+\int_{\eta_{2}}^{\pi}\right\} \frac{\operatorname{Sin}(n+k) u}{\tan \frac{1}{2} u} f^{\prime}(u) d u+ \\
& +\frac{2}{\pi} \int_{0}^{\pi} \operatorname{Cos}(n+k) u f(u) d u \\
& \quad=\frac{2}{\pi}\left\{J_{n}^{(1)}+J_{n}^{(2)}+J_{n}^{(3)}+J_{n}^{(1)}\right\}
\end{aligned}
$$

say, where $\eta_{1}=\frac{k_{1} \log (n+k)}{n+k}$ and $\eta_{2}=\frac{k_{1} \log n}{n}$.
It is plain that $J_{n}{ }^{(4)} \rightarrow 0$, as $n \rightarrow \propto$, uniformly with respect to $k$ and there. fore

$$
\begin{equation*}
\left\{\sum_{\nu=0}^{n}\left|J_{\nu}^{(4)}\right| q\right\}^{1 / q}=o\left(n^{3 / q}\right), \text { as } n \rightarrow \infty \tag{3-2}
\end{equation*}
$$

Now integrating by parts and using the condition (2.3) we get

$$
\begin{aligned}
& J_{n}^{(1)}=\sin (n+k) \eta_{1} \cot \frac{1}{2} \eta_{1} F\left(\eta_{1}\right)+o(1)- \\
& -\int_{0}^{\eta_{1}}\left\{(n+k) \cos (n+k) u \cot \frac{u}{2}-\frac{1}{2} \sin (n+k) u \operatorname{cosec}^{3} \frac{u}{2}\right\} F(u) d u \\
& = \\
& 0(1)+o(1) \cdot \int_{0}^{\eta_{1}}\left\{\frac{n+k}{u}+\frac{(n+k) u}{u^{2}}\right\} \cdot\left\{\frac{u}{\log \frac{1}{u}}\right\} d u
\end{aligned}
$$

$$
\begin{aligned}
& =o(1)+o(n+k) \int_{0}^{\eta_{1}} \frac{1}{\log \frac{T_{u}}{u}} d u \\
& =o(1), \text { as } n \rightarrow \propto, \text { unformly with respect to } k .
\end{aligned}
$$

Thus,
(3:3)

$$
\left\{\sum_{\nu=0}^{n}\left|J_{\nu}^{(1)}\right|^{q}\right\}^{1 / q}=0\left(n^{1 / q}\right), \text { as } n \rightarrow \infty,
$$

uniformly with respect to $k$.
Further using the condition (2.3), we have

$$
\begin{aligned}
J_{n}{ }^{(2)} & =\int_{\eta_{1}}^{\eta_{2}} f(u) \operatorname{Sin}(n+k) u \cot \frac{1}{\frac{1}{3}} u d u \\
& \leqslant \int_{\eta_{1}}^{\eta_{\eta_{2}}} \frac{f(u)}{u} d u \\
& =\left[\frac{F(u)}{u}\right]_{\eta_{1}}^{\eta_{2}}+\int_{\eta_{1}}^{\eta_{\eta_{1}}} \frac{F(u)}{u^{2}} d u \\
& =o(1), \text { as } n \rightarrow \infty
\end{aligned}
$$

uniformly w.r.t.k.
Therefore

$$
\left\{\sum_{\nu=0}^{n}\left|J_{\nu}^{(2)}\right| q\right\}^{1 / q}=o\left(n^{1 / q}\right), \text { as } n \rightarrow \infty
$$

uniformly w. r. t. k.
Now we denote by $C_{n}(\tau)$ the $n$-th Fourier sine coefficient of the odd function $x(t)$ which is equal to $f(t)$ in $(0, \tau)$ and zero in $(\tau, \pi)$, thus we have,

$$
\begin{aligned}
J_{n}^{(3)} & =\int_{\eta_{2}}^{\pi} \operatorname{Cot} \frac{1}{2} u\left\{\frac{d}{d t} \int^{u} f(t) \operatorname{Sin}(n+k) t d t\right\} d u \\
& =-\frac{1}{2} \pi \cot \frac{1}{2} \eta_{2} C_{n+k}\left(\eta_{2}\right)+\frac{\pi}{4} \int_{\eta_{2}}^{\pi} \operatorname{cosec}^{2} \frac{u}{2} C_{n+k}(u) d u .
\end{aligned}
$$

Using Minkowski's inequality, we write

$$
\begin{aligned}
\left\{\sum_{\nu=0}^{n}\left|J_{\nu}^{(s)}\right| q\right\}^{J / q} & \leqslant \cot \frac{1}{2} \eta_{2}\left\{\sum_{\nu=0}^{n}\left|C_{\nu+k}\left(\eta_{2}\right)\right|^{q}\right\}^{1 / q}+ \\
& +\int_{\eta_{2}}^{\pi} \operatorname{cosec}^{2} \frac{1}{2} u\left\{\sum_{\nu=0}^{n}\left|C_{\nu+k}(u)\right|^{q}\right\}^{1 / q} d u .
\end{aligned}
$$

Also by Hausdroff's inequality, we have

$$
\begin{aligned}
\left\{\sum_{\nu=0}^{n}\left|C_{\nu+k}(u)\right|^{q}\right\}^{1 / q} & \leqslant\left\{\int_{-\pi}^{\pi}|x(u)|^{p} d u\right\}^{1 / p} \\
& =\left\{\int_{-t}^{t}|f(u)|^{p} d u\right\}^{1 / p} \\
& =o(1) \cdot\left\{u^{1 / p} \cdot \log ^{1 / q} \frac{1}{u}\right\} .
\end{aligned}
$$

Hence

$$
\begin{align*}
\left\{\sum_{\nu=0}^{n}\left|J_{\nu}^{(3)}\right| q\right\}^{1 / q} & =O(1) \cdot\left\{\frac{1}{\eta_{2}}\left(\eta_{2}^{1 / p} \log ^{1 / q} \frac{1}{\eta_{2}}\right)\right\}+ \\
& +O(1) \cdot \int_{\eta_{2}}^{\pi} \bar{u}^{2} u^{1 / p} \log ^{1 / q} \frac{1}{u} d u \\
& =O(1) \cdot\left\{\eta_{2}{ }^{\frac{\Gamma}{p}-1} \log ^{1 / q} \frac{1}{\eta_{2}}\right\}+ \\
& +O(1) \int_{\eta_{2 ;} ; u^{\frac{1}{p}-2} \log ^{1 / q}}^{\pi} \frac{1}{u} d u \\
& =O(1) \cdot\left\{\eta_{\eta_{2}^{-1 / q}} \log ^{1 / q} \frac{1}{\eta_{2}}\right\} \\
& =o\left(n^{1 / q}\right), \text { as } n \rightarrow \infty
\end{align*}
$$

uniformly w.r.t. $k$.
Finally collection of $(3 \cdot 1),(3 \cdot 2),(3 \cdot 3),(3 \cdot 4)$ and (3.5) gives

$$
\left\{\sum_{\nu=0}^{n}\left|S_{\nu+k}\right| q\right\}^{3 / q}=o\left(n^{1 / q}\right), \text { as } n \rightarrow \infty,
$$

uniformly w. r. t. $k$.
Hence

$$
\sum_{\nu=0}^{n}\left|S_{\nu+k}\right| q=o(n), \text { as } n \rightarrow \propto,
$$

uniformly $w . r . t . k$.
This completes the proof of the theorem.

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## References

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$$
\left[\begin{array}{ll}
{[358}
\end{array}\right]
$$

## Some Theorems on Transform

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1. Introduction: In the present paper some theorems involving some functions and their transforms under any Fourier Kernel are obtained. The function $k(x)$ is said to be a Fourier Kernel if the following pair of reciprocal equations are simultaneously valid:

$$
\begin{align*}
g(x) & \left.=\int_{0}^{\infty} k(x)\right) f(y) d y  \tag{1}\\
f(x) & =\int_{0}^{\infty} k(x y) g(y) d y
\end{align*}
$$

To ensure the validity of (1) and (2), the set of conditions we use here consists of convergence conditions on $f(x)$ and $k(x)$ together with a functional equation satisfied by the Mellin transform of $k(x)$. $K(s)$ is said to be the Mellin transform of $k(x)$ if

$$
K(s)=\int_{0}^{\infty} k(x) x^{s-1} d x
$$

The functional equation satisfied by $K(s)$ is $K(s) K(1-s)=1$. [5]
The result given in Titchmarsh's Fourier Integrals and Mitra's result [4] follow as particular cases of the theorems of this paper.

The theorems have been formulated under rather stringent conditions. Some of the conditions may be relaxed and analytic continuation method may be applied.
2. If $\theta(x)$ and $\phi(x)$ are two Fourier Kernels and if the integral

$$
\int_{0}^{\infty} \phi(t) \theta(x t) d t \text { exists and is equal to } \phi_{\theta}(x), \text { say, then } \phi_{\theta}(x) \text { is also }
$$

a Fourier Kernel [5]
We shall require the following two throrems [5]:
I. Let $y^{c-1} f(y)$ belong to $L(0, \infty)$ and let $f(y)$ be continuous and of bounded variation in the neighbourhood of the point $y=x$. Let

$$
F(s)=\int_{0}^{\infty} f^{\prime}(x) x^{s-1} d x, s=c+i t \text {. Then }
$$

$$
f(s)=\frac{1}{2 \pi i} T \rightarrow \infty \int_{c-i}^{c+i T} F(s) x^{-s} d s
$$

II. Let $F(c+i u)$ belong to $L(-\infty, \infty)$ and let it be continuous and of bounded variation in the neighbourhood of the point $u=t$. Let

$$
\begin{aligned}
& f(x)=\frac{1}{2 \pi i} \int_{c-i \infty}^{0+i \infty} F(s) x^{-s} d x . \text { Then } \\
& F(s)=\lim _{T \rightarrow \infty} \int_{\frac{1}{T}}^{T} f(x) x^{s-1} d x \quad \text { for } s=c+i t
\end{aligned}
$$

In our theorems we shall assume that the conditions mentioned in these two theorems hold.
3. We shall adopt the following notations :

$$
\begin{aligned}
\phi_{\theta}(x) & =\int_{0}^{\infty} \phi(t) \theta(x t) d t \quad \text { and } \\
\left(\phi_{\theta}\right)(x) & =\int_{0}^{\infty} \theta(x t)_{\theta} \phi\left(\frac{1}{t}\right) \frac{d t}{t},
\end{aligned}
$$

where $\theta(x)$ and $\phi(x)$ are Fourier Kernels continuous in each arbitrary finite interval $(O, X)$ and the integrals converge uniformly in $(O, X)$, so that $\phi_{\theta}(x)$ exist and are continuous.

If $L(s)$ and $K(s)$ are the Mellin transforms of $\theta(x)$ and $\phi(x)$ respectively, then it is easy to see that $L(s) K(s)$ is the Mellin transform of $(00 \phi)(x)$ and $(\theta 0 \phi)(x)=$ $\left(\phi_{0 \theta}\right)(x)$ and $(\theta \circ \phi)(x)$ is also a Fourier Kernel.

We may remark here that if $\theta, \phi, \beta$ are three Fouricr Kernels then

$$
(\theta \circ(\phi \circ \beta))(x)=((00 \phi) \circ \beta)(\lambda)-(\theta 0 \phi \circ \beta)(x) .
$$

We may extend this result to $n$ Fourier kernels. Let $\theta_{1}, \theta_{2}, \ldots, \theta_{n}$ be $n$ Fourier kernels. Then $\left(\theta_{1} o \theta_{2} O \ldots o \theta_{n}\right)(x)$ is also a Fourier kernel. The $\theta^{\prime}$ 's may be permuted among themselves, and all $\theta$ 's are continuous in each arbitrary finite interval ( $0, X$ ) and the integrals converge uniformly. And

$$
\begin{gathered}
\quad\left(\theta_{1} \theta_{2} o \ldots o \theta_{n}\right)(x) \\
=\int_{0}^{\infty} \ldots \int_{0}^{\infty} \theta_{1}\left(t_{1}\right) \theta_{2}\left(t_{2}\right) \ldots t_{n-1}\left(t_{n-1}\right) \theta_{n}\left(\cdots \frac{x}{t_{1} t_{2} \ldots t_{n-1}}\right) \times \\
\times \frac{d t_{1} d t_{2} \ldots d t_{n-1}}{t_{1} t_{2} \ldots t_{n-1}} .
\end{gathered}
$$

If we take

$$
\begin{aligned}
& \theta_{1}(x)=\sqrt{x} J_{v_{1}}(x), \theta_{2}(x)=\sqrt{x} J_{v_{2}}(x), \ldots, \\
& \theta_{n}(x)=\sqrt{x} J_{v_{n}}(x), \text { then } \\
& \left(\theta_{1} 0 \theta_{2} 0 \ldots o \theta_{n}\right)(x)
\end{aligned}
$$

$$
\begin{gathered}
\left.=\sqrt{x} \int_{0} J_{v_{1}}\left(t_{1}\right) J_{v_{2}}\left(t_{2}\right) \ldots v_{v_{n-1} n-1}\right) J_{1_{n}}\left(\frac{x}{t_{1} t_{12} \ldots t_{n-1}}\right) \times \\
\times \frac{d t_{1} \ldots d t_{n-1}}{t_{1} \ldots t_{n-1}}
\end{gathered}
$$

$=\bar{\omega}_{1_{1}}, \ldots, \boldsymbol{v}_{n}(x)$, a Forier Kernel introduced by Bhatnagar [1]
In the theorems proved in the present paper, all the $\theta$ 's need not involve Bessel Functions of the first kind or sine or cosine functions. Recently Fox [3] has introduced a Kernel and our theorems remain true if we replace some or all the $\theta$ 's by Fox's Kernel provided the convergence of the integrals is secured and other conditions mentioned before hold.

We shall also adopt the notation :

$$
f_{\theta}(x)=\int_{0}^{\infty} f(t) \theta(x t) d t
$$

Provided $f(t)$ belongs to $L(0, \infty)$ and $\theta(x)$ is bounded and integrable in an arbitrary interval, $f_{\theta}(\lambda)$ always exists. We note that $f_{\theta}(x)$ is the Fourier Trans. form of $f(t)$ with regard to the Kernel $\theta(t)$.
4. We next prove a few theorems.

Theorem 1: Let $f(x)$ and $g(x)$ be two functions continuous and belonging to $L(0, \infty)$. Then

$$
\int_{0}^{\infty} f_{\theta}(x) g_{\theta}(x) d x=\int_{0}^{\infty} f(x) g(x) d x
$$

This is known theorem.
Theorem 2 : Let $f(x)$ and $g(x)$ be two continuous functions and $\theta(x)$, $\phi(x)$ iwo Fourier Kernels such that $\theta(x), \phi(x)$ are continuous in each arbitrary finite interval $(0, X)$ and the integral involving the Kernels i.e., $(\theta \circ \phi)(x)$ converge uniformly. Then

$$
\begin{aligned}
\int_{0}^{\infty} f_{\theta \circ \phi}(t) g_{\theta \circ \phi}(t) d t & =\int_{0}^{\infty} f_{\theta}(t) g_{\theta}(t) d t \\
& =\int_{0}^{\infty} f_{\phi}(t) g_{\phi}(t) d t \\
& =\int_{0}^{\infty} f(t) g(t) d t
\end{aligned}
$$

Before proving the theorem we prove the following lemma.
LEMMA : $\quad f_{\partial O \phi}(x)=\left(\frac{1}{x} f_{\theta}\left(\frac{1}{x}\right)\right)_{\phi}=\left(\frac{1}{x} f_{\phi}\left(\frac{1}{x}\right)\right)_{\theta}$.

Proof: We know that

$$
\begin{aligned}
f_{\theta 0 \phi}(x) & =\int_{0}^{\infty} f(t) d t \int_{0}^{\infty} \theta\left(\frac{1}{u}\right) p(x t u) \frac{d u}{u} \\
= & \int_{0}^{\infty} \theta\left(\frac{1}{u}\right) \frac{d u}{u} \int_{0}^{\infty} f(t) \phi(x t u) d t \\
& =\int_{0}^{\infty} \theta\left(\frac{x}{u}\right) \frac{d u}{u} \int_{0}^{\infty} f(t) \phi(t u) d t \\
& =\int_{0}^{\infty} \theta\left(\frac{x}{u}\right) f_{\phi}(u) \frac{d u}{u} \\
& =\int_{0}^{\infty} \theta(x u) f_{\phi}\left(\frac{1}{u}\right) \frac{d u}{u} \\
& =\left(\frac{1}{x} f_{\phi}\left(\frac{1}{x}\right)\right)_{\theta} .
\end{aligned}
$$

Similarly we can show that

$$
f_{\theta O \phi}(x)=\left(\frac{1}{x} f_{\theta}\left(\frac{1}{x}\right)\right)_{\phi}
$$

Thus the lemma is proved.
Proof of the theorem :

$$
\begin{aligned}
& \int_{0}^{\infty} f_{\theta \circ \phi}(x) g_{\theta o \phi}(x) d x \\
& =\int_{0}^{\infty}\left(\frac{1}{x} f_{\theta}\left(\frac{1}{x}\right)\right)_{\phi}\left(\frac{1}{x} g_{\theta}^{\prime}\left(\frac{1}{x}\right)\right)_{\phi} d x \\
& =\int_{0}^{\infty} \frac{1}{x} f_{\theta}\left(\frac{1}{x}\right) \frac{1}{x} g_{\theta}\left(\frac{1}{x}\right) d x \\
& =\int_{0}^{\infty} f_{\theta}(x) g_{\theta}(x) d x \\
& =\int_{0}^{\infty} f_{\phi}(x) g_{\phi}(x) d x \\
& =\int_{0}^{\infty} f(x) g(x) d x
\end{aligned}
$$

Corollary 1: Let $f(x)$ and $g(x)$ be two continuous functions belonging to $L(0, \infty)$ and $\theta_{1}(x), \theta_{y}(x), \ldots, \theta_{n}(x)$ be $n$ Fouricr Kernels such that $\theta_{r}(x), r=1,2, \ldots, n$, is continuous in each arbitrary finite interval ( $0, X$ ) and the integrals involving the Kernels i.e., $\left.\left(\theta_{1} 0\right)_{2} 0 \ldots o \theta_{r}\right)(x)$, converge uniformly. Then

$$
\begin{aligned}
& \int_{0}^{\infty} f_{\theta_{1}, 0 \theta_{2}, ~} \ldots o \theta_{n}(x) g_{\theta_{1}} o \theta_{2} o \ldots o \theta_{n}(x) d x \\
& =\int_{0}^{\infty} f_{\theta_{r}}(x) g_{\theta r}(x) d x,(r=\mathrm{I}, 2, \ldots, n) \\
& =\int_{0}^{\infty} f(x) g(x) d x
\end{aligned}
$$

Proof: Since the $\theta$ 's can be permuted among themselves

$$
\left(\theta_{1} 0 \theta_{2} 0 \ldots o \theta_{n}\right)(x)=\left(\theta_{r} 0\left(\theta_{1} 0 \theta_{2} 0 \ldots o \theta_{r-1} \circ \theta_{r+1} 0 \ldots o \theta_{n}\right)\right)(x) .
$$

Now ( $\left.\theta_{0} 0 \theta_{2} 0 \ldots \theta_{r-1} 0 \theta_{r+1} 0 \ldots o \theta_{r}\right)(x)$ is a Fourier Kernel continuous in each arbitrary finite interval ( $0, X$ ) (the integrals converging uniformly). Let us denote this Kernel by $\psi(x)$.
Therefore

$$
\left(\theta_{1} \circ \theta_{2} o \ldots o \theta_{n}\right)(x)=\left(\theta_{r} \circ \psi\right)(x) .
$$

Now

$$
\begin{aligned}
& \int_{0}^{\infty} f_{\theta_{1} o \theta_{1} o \ldots o \theta_{n}(x) g_{\theta 1} o \cdots o \theta_{n}(x) d x} \\
& =\int_{0}^{\infty} f_{\theta_{r} o \psi}(x) g_{\theta r} \circ \psi(x) d x \\
& =\int_{0}^{\infty} f_{\theta_{r}}(x) g_{\theta r}(x) d x \\
& =\int_{0}^{\infty} f(x) g(x) d x
\end{aligned}
$$

Coroliary 2: Under the conditions of corollary 1,

$$
\begin{aligned}
& \int_{0}^{\infty} f_{\theta_{1}} o \cdots o \theta_{n}(x) g_{\theta_{1}} o \cdots o \theta_{n}(x) d x \\
& =\int_{0}^{\infty} f_{\theta_{1}} o_{2} o \cdots o \theta_{r}(x) g_{\theta_{1}} o \theta_{2} o \cdots o \theta_{r}(x) d x \\
& =\int_{0}^{\infty} f_{\theta_{r+1}} o \theta_{r+2} o \cdots o \theta_{n}(x) g_{\theta_{r+1}} o \theta_{r+2} o \cdots o \theta_{n}(x) d x \\
& =\int_{0}^{\infty} f(x) g(x) d x, \text { where } r=1,2, \ldots, n-1 .
\end{aligned}
$$

Corollary 3 :

$$
\int_{0}^{\infty} f_{\theta(x)} g_{\theta 0 \phi}(x) d x=\int_{0}^{\infty} f\left(\frac{1}{t}\right) \frac{1}{l} g_{\phi}(t) d t
$$

provided $f(t)$ and $g(t)$ are continuous and belong to $L(0, \infty)$ and $\theta(x), \phi(x)$ are continuous in ( $0, X$ ) for all finite values of $X$ and the integral for ( $\theta \circ \phi)(x)$ cooverages uniformly, and the integral exist.

Proof: This can be proved by comparing the Mellin Transforms of both sides. We also get the following generalised result :

$$
\begin{aligned}
& \int_{0}^{\infty} f_{\theta_{1} 0 \theta_{1} o} \ldots o r_{r}^{(t)} g_{\theta_{0} \circ \theta_{2} o \ldots o \theta_{r} 0 \theta_{r+1} o \ldots o \theta_{n}}(t) d t \\
& =\int_{0}^{\infty} f\left(\frac{1}{t}\right) g_{r_{r 1} o \ldots o \theta_{n}}{ }^{(t) \frac{d t}{t}}
\end{aligned}
$$

under the conditions of the theorem.
In corollary 1 , putting $\left(\theta_{1} 0 \theta_{2} 0 \ldots 0 \theta_{n}\right)(x)=\bar{\omega}_{\nu_{1}}, \ldots, v_{n}(x)$ we get theorem 2 of Mitra [4]. In corollary 2, putting $r=1,2, \ldots, n-1$, in succession we get the corollary 1 of Mitra.

We note here that in all these theorems a change in the order of integration is permissible by De-la.Vallee Poussin's Theorem [2].
5. Let
(i) $\theta(x)$ and $\phi(x)$ be continuous functions of $x$ in ( $o, X)$ for all finite values of $X$ and let the integral involving $\theta$ and $\phi$, viz. $\phi_{\theta}(x)$ converge uniformly.
(ii) Let $f(x), g(x)$ and $h(x)$ be continuous and belong to $L(0, \infty)$.
(iii) $\phi_{\theta}(x)$ be bounded and integrable in $(0, \infty)$.
(iv) $f_{\theta}(x), g_{\phi}(x)$ and $h \phi_{\theta}(x)$ be absolutely integrable in ( $0, \infty$ ), where

$$
h_{\phi_{\theta}}(x)=\int_{0}^{\infty} h(t) \phi_{\theta}(x t) d t .
$$

Then
Theorem 3 :

$$
\begin{aligned}
& \int_{0}^{\infty} \int_{0}^{\infty}{ }^{h_{\phi}} \phi_{\theta}\left(x_{1}\right) f_{\theta}\left(x_{1} x_{2}\right) g_{\phi}\left(x_{2}\right) d x_{1} d x_{2} \\
& =\int_{0}^{\infty} \int_{0}^{\infty} h\left(t_{1}\right) f\left(t_{1} t_{2}\right) g\left(t_{2}\right) d t_{1} d t_{2}
\end{aligned}
$$

Proof: Now

$$
\begin{aligned}
& \int_{0}^{\infty} \int_{0}^{\infty} h_{\phi \theta}\left(x_{1}\right) f_{\theta}\left(x_{1} x_{2}\right) g_{\phi}\left(x_{2}\right) d x_{1} d x_{2} \\
& =\int_{0}^{h_{\phi}}\left(x_{1}\right) d x_{1} \int_{0}^{\infty} f_{\theta}\left(x_{1} x_{2}\right) g_{\phi}\left(x_{2}\right) d x_{2} \\
& =\int_{0}^{\infty} h_{\phi}\left(x_{1}\right) d x_{1} \int_{0}^{\infty}\left\{\int_{0}^{\infty} f\left(t_{1}\right) \theta\left(x_{1} x_{2} t_{1}\right)^{\prime} d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \phi\left(x_{2} t_{2}\right) d t_{2}\right\} d x_{2} \\
& =\int h_{\phi}\left(x_{1}\right) d x_{1} \int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) d t_{2} \int_{0}^{\infty} \theta\left(x_{1} x_{2} t_{1}\right) \phi\left(x_{2} t_{2}\right) d x_{2}
\end{aligned}
$$

$$
\left.\begin{array}{l}
=\int_{0}^{\infty} h_{\phi_{\theta}}\left(x_{1}\right) d x_{1} \int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \frac{d t_{2}}{t_{2}} \int_{0}^{\infty} \theta\left(\frac{x_{1} x_{2} t_{1}}{t_{2}}\right) \phi\left(x_{2}\right) d x_{2} \\
=\int_{0}^{\infty} h_{\phi_{\theta}}\left(x_{1}\right) d x_{1} \int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \phi_{\theta}\left(\frac{x_{1} t_{1}}{t_{2}}\right) \frac{d t_{2}}{t_{2}} \\
=\int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \frac{d t_{2}}{t_{2}} \int_{0}^{\infty} h_{\phi_{\theta}}\left(x_{1}\right) \phi_{\theta}\left(\frac{x_{1} t_{1}}{t_{2}}\right) d x_{1} \\
=\int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) h\left(\frac{t_{1}}{t_{2}}\right) d t_{t_{2}} \\
t_{t_{2}}
\end{array}\right]=\int_{0}^{\infty} g\left(t_{2}\right) d t_{2} \int_{0}^{\infty} f\left(t_{1} t_{2}\right) h\left(t_{1}\right) d t_{1} .
$$

The changes in the order of integrations are permissible by the conditions of the theorem.

Theorem 4: Under the same conditions

$$
\begin{aligned}
& \int_{0}^{\infty} \int_{0}^{\infty} h\left(t_{1}\right) f_{\theta}\left(t_{1} t_{2}\right) g_{\psi}\left(t_{2}\right) d t_{1} d t_{2} \\
& =\int_{0}^{\infty} \int_{0}^{\infty} h_{\phi}\left(t_{2}\right) f\left(t_{1} t_{2}\right) g\left(t_{2}\right) d t_{1} d t_{2}
\end{aligned}
$$

Proof: This can be proved in a similar manner.
Theorem 5 : Let
(i) $\theta(x)$ and $\phi(x)$ be continuous in $(0, X)$ for all finite values of $X$.
(ii) Let $f(x), g(x), \frac{1}{x} g(x)$ and $h(x)$ be continuous and belong to $L(0, \infty)$.
(iii) $\left(\begin{array}{l}\phi\end{array} 0\right)(x)$ be bounded and integrable in $(0, \infty)$.
(iv) $f_{\theta}(x), \frac{1}{x} G\left(\frac{1}{x}\right)$ and $h_{\phi_{0 \theta}}(x)$ be absolutely integrable in ( $0, \infty$ ), where

$$
G(x)=\int_{0}^{\infty} g(t) \phi\left(\frac{x}{t}\right) \frac{d t}{t}
$$

Then

$$
\begin{aligned}
& \int_{0}^{\infty} \int_{0}^{\infty} h_{\phi o_{\theta}}\left(x_{1}\right) f_{6}\left(x_{1} x_{2}\right) G\left(\frac{1}{x_{2}}\right) \frac{d x_{1} d x_{2}}{x_{2}} \\
& =\int_{0}^{\infty} \int_{0}^{\infty} h\left(t_{1}\right) f\left(t_{1} t_{2}\right) g\left(t_{2}\right) d t_{1} d t_{2} \\
& {[365]}
\end{aligned}
$$

Proof :

$$
\begin{aligned}
& \int_{0}^{\infty} \int_{0}^{\infty} h_{\phi O \theta}\left(x_{1}\right) f_{\theta}\left(x_{1} x_{2}\right) G\left(\frac{1}{x_{2}}\right) \frac{d x_{1} d x_{2}}{x_{2}} \\
& =\int_{0}^{\infty} h_{\phi o \theta}\left(x_{1}\right) d x_{1} \int_{0}^{\infty}\left\{\int_{0}^{\infty} f\left(t_{1}\right) \theta\left(x_{1} x_{2} t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \phi\left(\frac{1}{x_{2} t_{2}}\right) \frac{d t_{2}}{t_{2}}\right\} \frac{d x_{2}}{x_{2}} \\
& \quad=\int_{0}^{\infty} h_{\phi o \theta}\left(x_{1}\right) d x_{1} \int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \frac{d t_{2}}{t_{2}} \int_{0}^{\infty} \theta\left(x_{1} x_{2} t_{1}\right) \phi\left(\frac{1}{x_{2} t_{2}}\right) \frac{d x_{2}}{x_{2}} \\
& \quad=\int_{0}^{\infty} h_{\phi o \theta}\left(x_{1}\right) d x^{1} \int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \frac{d t_{2}}{t_{2}} \int_{0}^{\infty} \theta\left(\frac{x_{1} x_{2} t_{1}}{t_{2}}\right) \phi\left(\frac{1}{x_{2}}\right) \frac{d x_{2}}{x_{2}} \\
& \quad=\int_{0}^{\infty} h_{\phi_{0 \theta}}\left(x_{1}\right) d x_{1} \int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \frac{1}{t_{2}}(\phi o \theta)\left(\frac{x_{1} t_{1}}{t_{2}}\right) d t_{2} \\
& =\int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \frac{1}{t_{2}} d t_{2} \int_{0}^{\infty} h_{\phi o \theta}\left(x_{1}\right)(\phi o \theta)\left(\frac{x_{1} t_{1}}{t_{2}}\right) d x_{1} \\
& =\int_{0}^{\infty} f\left(t_{1}\right) d t_{1} \int_{0}^{\infty} g\left(t_{2}\right) \frac{1}{t_{2}} h\left(\frac{t_{1}}{t_{2}}\right) d t_{2} \\
& =\int_{0}^{\infty} h\left(t_{1}\right) f\left(t_{1} t_{2}\right) g\left(t_{2}\right) d t_{1} d t_{2} .
\end{aligned}
$$

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# Convolution in a General Transform 

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## Abstract

In this paper Convolution of two functions have been expressed as an integral of transforms of these functions with respect to a pair of inverse transforms. A result of Mellin transform in terms of fractional integral is illustrated by examples.

## Introduction

1. Convolution of two functions $f(t)$ and $g(t)$ is defined by Erde'lyi [4, page 8]

$$
f * g=\int_{0}^{t} f(u) g(t-u) d u \quad t \geqslant 0
$$

for the interval $(o, t)$
Ruel. V. Churchill [3, page 274] has defined it over the interval ( $-\pi, \pi$ ). Convolution or Faltung of the functions $f$ and $g$ over the interval $(-\infty, \infty)$ is expressed as Fourier Transforms of functions $f$ and $g$ [5, page 24] as

$$
f * g=\int_{-\infty}^{\infty} f(x) g(u-x) d x=\int_{-\infty}^{\infty} F(t) G(t) e^{-i u t} d t
$$

Similar results for intervals $(0, x)$ and $(x, \infty)$ are known [5, page 31] for Laplace transforms.

Convolution of two functions over different intervals have been expressed in terms of integrals involving transforms of the functions $f$ and $g$.

The object of this paper is to express the Convolution of two functions $f$ and $g$ as an integral involving transforms of these functions. This will be done by establishing a theorem. Some particular cases will give known results while others are believed to be new. The theorem is illustrated by two examples. The approach is formal.

## 2. Theorem If

(i) $k_{1}(x)$ and $k_{2}(x)$ are kernels of a transform and its inverse, having $C_{,}$and $C_{2}$ as the limits of integration so that

$$
\begin{align*}
F(y) \equiv K_{1}\{f(x) ; y\} & =\int_{c_{1}} f(x) k_{1}(x y) d x \\
f(x) \equiv K_{2}\{F(y) ; x\} & =\int_{c_{2}} F(y) k_{12}(x y) d y
\end{align*}
$$

(ii) $f(t) g(x-t)$ is integrable over the interval $C_{1}$
(iii) the integral

$$
\int_{c_{2}} \kappa_{1}\{g(\xi) ; y\} k_{2}(y x) d y
$$

is absolutely convergent
then
the convolution of two functions $f(t)$ and $g(t)$ over the interval $C_{1}$ is given by

$$
\int_{c_{1}} f(t) g(x-t) d t=\int_{c_{2}} G(y) r(x, y) d y
$$

where

$$
G(y)=K_{1}\{g(\xi) ; y\} \text { and } r(x, y)=\int_{c_{1}} f(t) k_{2}\{(x-t) y\} d t
$$

Proof: Because of (i) we can write

$$
g(x-t)=\int_{c_{2}} K_{1}\{g(\xi) ; y\} K_{2}\{(x-t) y\} d y
$$

Using (2.5) in L. H. S. of $(2 \cdot 3)$ we get

$$
\int_{c_{1}} f(t) g(x-t) d t=\int_{c_{2}} f(t)\left[\int_{c_{2}} K_{1}\{g(\xi) ; y\} k_{2}\{(x-t) y\} d y\right] d t
$$

Changing the order of integration which is permissible under the conditions of the theorem we have

$$
\int_{c_{1}} f(t) g(x-t) d t=\int_{c_{2}} K_{1}\{g(\xi) ; y\}\left[\int_{c_{1}} f(t) K_{2}\{(x-t) y\} d t\right] d y
$$

Using (2•4) we arrive at (2•3).
3. Particular Cases: In cases where $k_{2}\{(x-t) y\}$ can be broken up into factors of the type $k_{2}(x y) k_{1}(y t)$, the function $r(x, y)$ is further simplified. This can be seen in the cases of Fourier transform, Laplace transform and Mellin Inverse transform. Other cases of Mellin transform, Fourier Cosine, Fourier Sine are believed to be new.
(1) Fourier Transform. Let $k_{1}(x)=e^{i x}, k_{2}(x)=e^{-i x}, C_{1}=C_{2}=(-\infty, \infty)$ then $r(x, y)$ of (2.4) becomes $e^{-i x y} F(y)$ and (2•3) reduces to the known result [5, page 24].

$$
\int_{-\infty}^{\infty} f(t) g(x-t) d t=\int_{-\infty}^{\infty} G(y) F(y) e^{-i x y} d y
$$

Interchanging $k_{1}(x)$ and $k_{2}(x)$ we have

$$
\int_{-\infty}^{\infty} G(x-t) F(t) d t=\int_{-\infty}^{\infty} f(y) g(y) e^{i x y} d y
$$

(ii) Laplace Transform. Let $k_{1}(x)=e^{-p x}, k_{2}(p)=\frac{e^{p x x}}{2 \pi i} C_{1}$ is $(0, \infty)$ and $C_{2}$ is $(c-i \infty$ to $c+i \infty$ ) then $r(x, y)=e^{x y y} F(y)$ and (2.3) gives known result [5, page 31]. Interchanging $k_{1}(x), C_{1}$ and $k_{2}(\lambda), C_{2}$ in (2:3) we have

$$
\frac{1}{2 \pi i} \int_{c-i \infty}^{c+i \infty} G(x-t) F(t) d t=\int_{0}^{\infty} g(y) f(y) e^{-x y} d y
$$

(iii) Mellin Inverse. Let $k_{1}(s x)=x^{-s} / 2 \pi i, k_{2}(s x)=x^{8-1} C_{1}$ is $(c-i \infty$ to $c+i \infty)$, $C_{2}$ is $(o$ to $\infty)$, then $r(x, y)=y^{2 v-1} f(y)$ and (2.3) assumes a known form [6, page 527.
(iv) Mellin Transform. Let $k_{1}(x s)=x^{s-1}, k_{2}(s x)=x^{-s} / 2 \pi i$

Now
(3.1)

$$
g(x-t)=\frac{1}{2 n \pi i} \int_{c-i \infty}^{0+i_{\infty}} G(y)(x-t)^{-y} d y
$$

and this exists for $\quad t<\boldsymbol{x}$
Modyfying accordingly

$$
r(x, y)=\Gamma(1-y) I_{x}^{(1-y)} f(t) . \quad 0<y<1
$$

where $I_{x}^{(1-y)} f$ denotes the operator of Reimann-Liouville fractional integral operator of order $(1-y)$ on $f(t)$.
Ultimately ( $2 \cdot 3$ ) becomes

$$
\begin{array}{r}
\int_{0}^{x} f(t) g(x-t) d t=\frac{1}{2 \pi}-\int_{c-i \infty}^{c+i \infty} G(y) \Gamma(1-y) I_{x}^{(1-y)} f(t) d y . \\
0<y<1
\end{array}
$$

Above result can be easily verified by taking

$$
f(t)=t^{\prime-1} \text { and } g(t)=e^{-t}, \operatorname{Re} v>0 .
$$

(v) Fourier Cosine Transform. Let $k_{1}(x)=k_{2}(x)=\sqrt{\frac{2}{\pi}} \operatorname{Cos} x, F_{c}(y)$ and $G_{c}(y)$ denote the Fourier Cosine transforms of $f(t)$ and $g(t)$ and let

$$
F_{s, x}(y)=\int_{0}^{x} f(t) \sin (t y) d t
$$

and

$$
F_{c, x}(y)=\int_{0}^{x} f(t) \operatorname{Cos}(t y) d t
$$

then

$$
r(x, y)=\sqrt{\frac{2}{\pi}} \operatorname{Cos}(x y) F_{c, x}(y)+\sqrt{\frac{\overline{2}}{\pi}} \operatorname{Sin}(x y) F_{s, x}(y)
$$

Thus theorem becomes

$$
\begin{align*}
\int_{0}^{x} f(t) g(x-t) d t & =\sqrt{\frac{2}{\pi}} \int_{0}^{\infty} G_{c}(y) F_{c, x}(y) \operatorname{Cos}(x y) d y \\
& +\sqrt{\frac{\pi}{\pi}} \int_{0}^{\infty} G_{c}(y) F_{s, x}(y) \operatorname{Sin}(x y) d y \\
& {[369] }
\end{align*}
$$

(vi) Similar result for Fourier Sine transform will be

$$
\begin{gather*}
\int_{0}^{x} f(t) g(x-t) d t=\sqrt{\frac{2}{\pi}} \int_{0}^{\infty} G_{s}(y) F_{c, x}(y) \operatorname{Sin}(x y) d y \\
-\sqrt{\frac{2}{\pi}} \int_{0}^{\infty} G_{s}(y) F_{s, x}(y) \operatorname{Cos}(x y) d y
\end{gather*}
$$

Above results can be verified by taking functions

$$
f(t)=e^{-\alpha t}, g(t)=e^{-\beta t}, \alpha>\beta>0 .
$$

4. Example 1: Let $f(t)=\operatorname{Cos}(a t), g(t)=e^{-\beta t}$

Usirg the result [2, page 189 V ol. II].
$(4 \cdot 1) I_{x^{(1-y)}} \operatorname{Cos}(a t)=\frac{x^{1-y}}{2 \Gamma(2-y)}\left[{ }_{1} F_{1}(1 ; 2-y ; i a x)+{ }_{1} F_{1}(1 ; 2-y ;-i a x)\right]$
$1-y>0$.
Using [1, page 312, Vol. I]
Mellin transform of $e^{-\beta t}$ is

$$
G(y)=\frac{\Gamma(y)}{\beta y} y>0 .
$$

Substituting these values in (3.2) we have

$$
\begin{aligned}
& \int_{0}^{x} \operatorname{Cos}(a t) e^{-\beta(x-t)} d t=\frac{1}{2 \pi i} \int_{c-i \infty}^{c+i \infty} \frac{\Gamma(y) \Gamma(1-y) x^{1-y}}{\beta y 21(2-y)} \\
& \quad \times\left[{ }_{1} F_{1}(1 ; 2-y ; i a x)+{ }_{1} F_{1}(1 ; 2-y ;-i a x)\right] d y
\end{aligned}
$$

Evaluating the L. H. S. and simplifying the R. H. S. we get

$$
\begin{array}{r}
\frac{1}{2 \pi i} \int_{c-i_{\infty}}^{c+i_{\infty}} \frac{\Gamma(y)}{(1-y)}\left[{ }_{1} F_{1}(1 ; 2-y ; i a x)+{ }_{1} F_{1}(1 ; 2-y ;-i a x)\right] x^{-y} \beta^{-y} d y \\
=\frac{2}{x\left(a^{2}+\beta^{2}\right)}\left\{\beta \operatorname{Cos}(a x)+a \operatorname{Sin}(a x)-\beta e^{-\beta x}\right\} 0<y<1 \\
|\beta|<|a|
\end{array}
$$

Replacing $\beta$ by $z, x$ by $\lambda$, we have

$$
\begin{array}{r}
\frac{1}{2 \pi i} \int_{c-i \infty}^{c+i i_{\infty}} \frac{\Gamma(y)}{(1-y)}\left[{ }_{1} F_{1}(1 ; 2-y ; i a \lambda)+{ }_{1} F_{1}(1 ; 2-y ;-i a \lambda)\right] \lambda^{-y} z^{-y} d y \\
=\frac{2}{\lambda\left(a^{2}+z^{2}\right)}\left\{z \operatorname{Cos}(a \lambda)+a \operatorname{Sin}(a \lambda)-z e^{-\lambda z}\right\} 0<\operatorname{Re} y<1 \\
|z|<|a|
\end{array}
$$

This result can be verified by evaluating the integral with the help of the residue theorem, the contour being a line parallel to imaginary axis, cutting the real axis at $c$ such that $o<c<1$, encircling the poles viz. $(0,-1,-2 \ldots \ldots$ )
Example 2: Let $f(t)=e^{\alpha t}, g(t)=\operatorname{Sin} \beta t$, then using [2, page 187].

$$
I_{x^{(1-y)}} f=\frac{x^{1-y}}{\Gamma}(2-y){ }_{1} F_{1}\{1 ; 2-y ; \alpha x\} \quad 1-y>0
$$

also using [1, page 317] Mellin transform of $g(t)$, becomes

$$
G(y)=\frac{\Gamma(y)}{\beta^{y}} \operatorname{Sin}\left(\frac{\pi}{2} y\right)-1<\operatorname{Re} y<1
$$

Substituting these values in (3.2) and simplifying we have

$$
\begin{aligned}
& \frac{1}{2 \pi i} \int_{c-i \infty}^{c+i_{\infty}} \frac{\Gamma(y)}{(1-y)} \operatorname{Sin}\left(\frac{\pi}{2} y\right){ }_{1} F_{\perp}\{1 ; 2-y ; a x\} x^{-y} \beta^{-y} d y \\
& =\frac{\beta\left(e^{\alpha x}-\operatorname{Cos} \beta x\right)-a \operatorname{Sin} \beta x}{x\left(\alpha^{2}+\beta^{2}\right)} 0<\operatorname{Re} y<1,|\beta|<|\alpha|
\end{aligned}
$$

Replacing $\beta$ by $z, x$ by $\lambda$, we have

$$
\begin{aligned}
& \frac{1}{2 \pi i} \int_{c-i \infty}^{c+i \infty} \frac{\Gamma(y)}{(1-y)} \operatorname{Sin}\left(\frac{\pi}{2} y\right){ }_{1} F_{1}\{1 ; 2 y ; \alpha \lambda\} \lambda^{-y} z^{-y} d y \\
& =\frac{z\left(e^{\alpha \lambda}-\operatorname{Cos} \lambda z\right)-\alpha \operatorname{Sin} \lambda z}{\lambda\left(\alpha^{2}+z^{2}\right)} 0<\operatorname{Re} y<1,|z|<|\alpha|
\end{aligned}
$$

This example can also be verified as in Ex. 1.

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# Superposition of Perfect Fluid and Electromagnetic Fields in Sperically Symmetric Space-Time in General Relativity 

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#### Abstract

In obtaining solutions of the field equations of general relativity the energymomentum tensor employed is either that of material distribution or of electromagnetic distribution. In this paper we examine whether the superposition of the two fields is consistent within the frame-work of general relativity. It has been shown that such a superposition is possible in case when the material field consists of perfect fluid distribution and the space-time is spherically symmetric. In the special case of the distribution of discrete particles the density and the electromagnetic energy have been expressed in terms of the metric potentials.


## Introduction

The form of the material energy tensor in general relativity in case of perfect fluid distribution is given by

$$
\mathcal{T}_{j}^{i}=(\varepsilon+p) u^{i} u_{j}+\delta_{j}^{i} p, \quad u^{i} u_{i}=-1,
$$

where $\varepsilon, p, u^{i}$ are respectively the matter density, pressure and the unit flow vector. On one hand when $p=0$ it leads to a distribution of discrete particles and on the other, when $T=\mathcal{T}_{i}{ }^{i}=0$ it gives the photon fluid distribution corresponding to the maximum pressure-density ratio, viz. $p / \varepsilon=\frac{1}{3}$. Other cases of perfect fluid distributions having a given equation of state lie between these two extreme cases. The form of the material-energy tensor in case of viscous fluids has recently been given by Lichnerowicz (1955). Also we know that the form of the energy tensor for the electromagnetic distribution is given by

$$
E_{j}^{i}=-F_{j a} F^{i a}+\frac{1}{4} S_{j}^{i} F_{a b} F^{\alpha b},
$$

where $F_{i j}$ is the skew-symmetric electromagnetic field tensor.
The various solutions in general relativity corresponding to non-empty fields have been obtained by using either of these field equations. But the form of the energy tensor when the matter and the electromagnetic field exist concurrently has not been specified in general relativity. In solving a particular problem of this nature a linear superposition of the two fields $T_{j}{ }^{i}$ and $E_{j}{ }^{i}$ has iecently been considered by Shah and Vaidya (1967). In this paper we examine a more ger eral case of superposition of the perfect fluid distribution and the electromagnetic field within the frame-work of general relativity. For this purpose we take the total
energy tensor as the linear sum of the two tensors given in ( $1 \cdot 1$ ) and (1-2). The case of discrete particle distribution has also been examined. It is shown in this paper that the field equations thus obtained are consistent. The pressure and density are modified by the presence of the electromagnetic field. Nordstrom solution (2) is given as a particular case as is expected.

## Metric and field equations

We consider the spherically symmetric space-time given by

$$
d s^{2}=-e^{\alpha} d t^{2}+e^{\beta} d r^{2}+r^{2}\left(d g^{2}+\sin ^{2} \theta d p^{2}\right) .
$$

where $a, \beta$ are functions of $r$ and $t$. The relativistic field equations are

$$
R_{j}^{i}-\frac{1}{2} R \delta_{j}^{i}=-8 \pi T_{j}^{i}
$$

where we take

$$
T_{j}{ }^{i}=T_{j}{ }^{i}+E_{j}{ }^{i} .
$$

For the metric ( $2 \cdot 1$ ) the field equations $(2 \cdot 2 a)$ reduce to

$$
\begin{align*}
& -8 \pi\left[(\varepsilon+p) u^{1} u_{1}+p+E_{1}^{1}\right]=\frac{1}{r_{2}}-e^{-\beta}\left(\frac{1}{r^{2}}+\frac{a_{1}}{r}\right), \\
& -8 \pi\left[(\varepsilon+p) u^{2} u_{2}+p+E_{2}^{2}\right]=-8 \pi\left[(\varepsilon+p) u^{3} u_{3}+p+E_{3}^{3}\right] \\
& =e^{-\alpha}\left(\frac{\beta_{41}}{2}+\frac{\beta_{4}{ }^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}{ }^{2}}{4}-\frac{\alpha_{11}}{2}\right) \\
& -8 \pi\left[(\varepsilon+p) u^{4} u_{4}+p+E_{4}^{4}\right]=\frac{1}{r^{2}}+e^{-\beta}\left(\frac{\beta_{1}}{r}-\frac{1}{r_{2}}\right) \text {, } \\
& -8 \pi\left[(\varepsilon+p) u^{1} u_{4}+E_{1}^{1}\right] e \beta=8 \pi e^{\alpha}\left[(\varepsilon+p) u^{4} u_{1}+E_{1}^{4}\right]=-\frac{\beta_{4}}{r} \\
& 8 \pi E_{2}^{1}=8 \pi e^{-\beta}\left[\frac{F_{23}}{r^{2}} F_{13} \operatorname{cosec}^{2} \theta-e^{-\alpha} F_{94} F_{14}\right]=-8 \pi(\varepsilon+p) u^{1} u_{2}, \\
& 8 \pi E_{3}^{1}=8 \pi e^{-\beta}\left[\frac{F_{12}}{r^{2}} F_{32}-e^{-a} F_{14} F_{34}\right] \quad=-8 \pi(\varepsilon+p) u^{1} u_{3}, \\
& 8 \pi E_{3}^{2}=\frac{8 \pi}{r^{2}}\left[e^{-\beta} F_{21} F_{31}-e^{-\alpha} F_{24} F_{34}\right] \quad=-8 \pi(\varepsilon+p) u^{2} u_{3}, \\
& 8 \pi E_{4}^{2}=\frac{8 \pi}{r^{2}}\left[e^{-\beta} F_{21} F_{41}+\underset{r^{2}}{F_{23}} F_{43} \operatorname{cosec}^{2} \theta\right] \quad=-8 \pi(\varepsilon+p) u^{2} u_{4}, \\
& 8 \pi E_{4}^{s}=\frac{8 \pi}{r^{2}}\left[\left.e^{-\beta} F_{31} F_{41}+\frac{F_{32}}{r^{2}} F_{42} \right\rvert\, \operatorname{cosec}^{2} \theta=-8 \pi(\varepsilon+p) u^{3} u_{4},\right.
\end{align*}
$$

where

$$
a_{1}=\frac{\partial r}{\partial r}, \alpha_{4}=\frac{\partial r}{\partial l}, \text { etc. }
$$

From equations $(2 \cdot 3 a, \ldots, i)$ we have

$$
\begin{array}{r}
E_{3}^{3}-E_{2}^{2}=(\varepsilon+f)\left(u^{2} u_{2}-u^{3} u_{3}\right) \\
{\left[\begin{array}{c}
373
\end{array}\right]}
\end{array}
$$

$$
\begin{align*}
& =\frac{1}{r^{\overline{2}}}\left[e^{-\beta}\left(F_{13}\right)^{2}-e^{-a}\left(F_{34}\right)^{2}\right] \operatorname{cosec}^{2} \theta-\frac{1}{r^{2}}\left[e^{-\beta}\left(F_{12}\right)^{2}-e^{-\alpha}\left(F_{24}\right)^{2}\right], \\
& \left(u^{1}\right)^{2}=(s+p)^{-1} e^{2} \beta k, \\
& \text { where } \\
& k=\frac{1}{r^{2}} \frac{\left[\left(E_{2}^{1}\right)^{2}-\left(E_{\frac{1}{1}}^{1}\right)^{2} \operatorname{cosec}^{2} \theta\right]}{E_{3}^{\frac{1}{3}}-E_{2}^{2}} \\
& \left(u_{2}\right)^{2}=(\varepsilon+p)^{-1} \frac{\left(E_{2}^{2}\right)^{2}}{k}, \\
& \left(u_{i}\right)^{2}=(\varepsilon+p)^{-1} \frac{\left(E_{3}^{1}\right)^{2}}{k}, \\
& \left(u_{4}\right)^{2}=(\varepsilon+p)^{-1} r^{4} \frac{\left(E_{4}^{2}\right)^{2}}{\left(E_{2}^{2}\right)^{2}} k, \\
& -p=e \beta k+E_{1}^{1}+\frac{1}{8 \pi r^{2}}-\frac{e^{-\beta}}{8 \pi}\left(\frac{1}{r^{2}}+\frac{\alpha_{1}}{r}\right), \\
& \left.-\varepsilon=\frac{1}{r^{2} k}\left[E_{2}^{1}\right)^{2}+\left(E_{8}^{1}\right)^{2} \operatorname{cosec}^{2} 6\right]-e^{-\alpha} r^{4} \frac{\left(E_{4}^{2}\right)^{2}}{\left(E_{2}^{1}\right)^{2}} k-E_{1}^{1}-\frac{1}{8 \pi r^{2}}+\frac{e^{-\beta}}{8 \pi}\left(\frac{1}{r^{2}}+\frac{\sigma_{1}}{r}\right)(2 \cdot 6 b) \\
& \left(E_{\mathrm{g}}^{1} E_{4}^{2}\right)^{2}=\left(E_{\frac{1}{2}}^{1} E_{4}^{3}\right)^{2} \sin ^{4} \theta, \\
& \left(E_{2}^{1} E_{3}^{1}\right)^{2}=r^{4}\left(E_{8}^{2}\right)^{2} k^{2}, \\
& 8 \pi\left[E_{1}^{1}-E_{2}^{2}+e \beta k-\frac{\left(E_{2}^{1}\right)^{2}}{r^{2} k}\right]=e^{-\alpha}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right) \\
& +e^{-\beta}\left(\frac{\beta_{1}+\alpha_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}^{2}}{4}-\frac{\alpha_{11}}{2}+\frac{1}{r^{2}}\right)-\frac{1}{r^{2}}  \tag{27c}\\
& 8 \pi\left[E_{1}^{1}-E_{4}^{4}+e^{\beta} k+e^{-a} \frac{\left(E_{ \pm}^{2}\right)^{2}}{\left(\overline{E_{2}^{3}}\right)^{2}} r^{1} k\right]=\frac{e^{-\beta}}{r}\left(\alpha_{1}+\beta_{1}\right)-\frac{1}{r^{2}}, \\
& 8 \pi r^{2} \frac{E_{4}^{2}}{E_{2}^{2}} k= \pm\left(8 \pi E_{4}^{1}-e^{-\beta} \frac{\beta_{4}}{r}\right) \tag{e}
\end{align*}
$$

Thus we have five algebraic equations, viz. $(2 \cdot 7 a, b, c, d, e)$ in eight unknown, viz . $F_{i j}, a, \beta$. Apart from these algebraic conditions $F_{i j}$ 's have to satisfy EinsteinMaxwell differential equations

$$
F_{[i j, k]}=0
$$

and

$$
F_{; j}^{i j}=J^{i},
$$

where a suffix preceded by a comma (,) stands for partial differentiation and that by a semi-colan (; ) indicates covariant differentiation.

The number of algebraic equations is less than the number of unknowns. In the next section we shall investigate the consequences of imposing restrictions on the unknowns.

## Radial flow of Matter

We assume that the flow of matter takes place along the radial direction only, i.e.

$$
u^{2}=u^{3}=0
$$

Equations $(1 \cdot 1),(2 \cdot 3 b, e, f . g, h, i)$ and $(2 \cdot 4)$ reduce to

$$
\begin{align*}
& u^{1} u_{1}+u^{4} u_{4}=-1, \\
& -8 \pi\left(p+E_{2}^{2}\right)=-8 \pi\left(p+E_{3}^{3}\right) \\
& =e^{-a}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\left.\beta_{1}-\frac{\alpha_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}^{2}}{4}-\frac{a_{11}}{2}\right)}{2} \begin{array}{l}
\frac{F_{23}}{r^{2}} F_{13} \operatorname{cosec}^{2} \theta-e^{-\alpha} F_{14} F_{24}
\end{array}=0,\right. \\
& \frac{F_{12}}{r^{2}} F_{32}-e^{-\alpha} F_{14} F_{34} \quad=0, \\
& e^{-\beta} F_{21} F_{31}-e^{-\alpha} F_{24} F_{34} \quad=0, \\
& e^{-\beta} F_{21} F_{41}+\frac{F_{23}}{r^{2}} F_{43} \operatorname{cosec}^{2} \theta=0, \\
& e^{-\beta} F_{31} F_{41}+\frac{F_{32}}{r^{2}} F_{42} \quad=0, \\
& e^{-\beta\left(F_{13}\right)^{2}-e^{-\alpha}\left(F_{34}\right)^{2}-\left[e^{-\beta}\left(F_{12}\right)^{2}-e^{-a}\left(F_{24}\right)^{2}\right] \sin ^{2} \theta=0 .}
\end{align*}
$$

From (3.2d) and ( $3 \cdot 2 g$ ) we get

$$
\left.\begin{array}{l}
\left(F_{13}\right)^{2}=-e e^{\beta-a}\left(F_{24}\right)^{2} \sin ^{2} \theta, \\
\left(F_{12}\right)^{2}=-e^{\beta-a}\left(F_{34}\right)^{2} \sin ^{2} \theta,
\end{array}\right\} \quad \text { Or } \quad\left\{\begin{array}{l}
\left(F_{13}\right)^{2}=e^{\beta-a}\left(F_{34}\right)^{2}, \\
\left(F_{12}\right)^{2}=e \beta-\alpha\left(F_{24}\right)^{2}
\end{array}\right.
$$

Since the first set of equations (3.2h) gives imaginary values for $F_{3_{1}}$ and $F_{21}$ we take the second set of $(3 \cdot 2 h)$, i.e.

$$
\begin{align*}
F_{12} & = \pm e^{\frac{\beta-\alpha}{2}} F_{14} \\
F_{13} & = \pm e^{\frac{\beta-a}{2}} F_{34} .
\end{align*}
$$

In view of $(3 \cdot 3 a, b)$ the equations $(3 \cdot 2 b),(3 \cdot 2 c),(3 \cdot 2 e)$ and $(3 \cdot 2 f)$ reduce to

$$
\begin{align*}
& \frac{\left(F_{23} F_{34}\right)^{2}}{r^{4}} \operatorname{cosec}^{4} \theta-e^{-a-\beta}\left(F_{14} F_{24}\right)^{2}=0, \\
& \left(F_{23} F_{24}\right)^{2} r^{4}-e^{-a-\beta}\left(F_{14} F_{34}\right)^{2}=0,
\end{align*}
$$

and ( $3.2 d$ ) reduces to an identity.
Now equations $(3 \cdot 4 a, b)$ lead to the following cases :
Case (ia)

$$
F_{14}, F_{23}, F_{24}, F_{34} \neq 0
$$

In this case $(3 \cdot 3 a)$ and $(3 \cdot 3 b)$ give

$$
F_{34}= \pm F_{24} \sin \theta,
$$

[ 375 ]

$$
F_{23}= \pm r^{2} e^{\frac{-\alpha-\beta}{2}} F_{14} \sin \theta,
$$

and further we get

$$
\left.\begin{array}{l}
E_{\perp}^{\mathrm{i}}=e^{-\alpha-\beta}\left(F_{14}\right)^{2}-\frac{2}{r^{2}} e^{-\alpha}\left(F_{24}\right)^{2},  \tag{3•6}\\
E_{2}^{2}=E_{\mathrm{s}}^{\mathrm{a}}=-e^{-a-\beta}\left(F_{14}\right)^{2}, \\
E_{4}^{4}=e^{-\alpha-\beta}\left(F_{14}\right)^{2}+\frac{2}{r^{2}} e^{-\alpha}\left(F_{24}\right)^{2} \\
E_{4}^{1}= \pm \frac{2}{r^{2}} e^{\frac{-\alpha-\beta}{2}}\left(F_{24}\right)^{2}
\end{array}\right\}
$$

From (3.5a, b), (3.6), (2.3a, c, d), (3•1) and (3.2a) we get

$$
\begin{gather*}
-8 \pi p=-8 \pi e^{-\alpha-\beta\left(F_{14}\right)^{2}+e^{-\alpha}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}{ }^{2}}{4}-\frac{\alpha_{13}}{2}\right)( } \begin{array}{c}
8 \pi \varepsilon=3(8 \pi) e^{-\alpha-\beta}\left(F_{14}\right)^{2}+\frac{2}{r^{2}}-e^{-\alpha}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+ \\
+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\mu_{1}^{2}-\alpha_{1} \beta_{1}}{4}+\frac{\alpha_{11}}{2}-\frac{2}{r^{2}}\right)
\end{array},
\end{gather*}
$$

$$
\begin{equation*}
\left(u^{1}\right)^{2}=\frac{(\varepsilon+p)^{-1}}{8 \pi} e^{-\beta}\left[-8 \pi p+e^{-\beta}\left(\frac{1}{r^{2}}+\frac{\alpha_{1}}{r}\right)-\frac{1}{r^{2}}+8 \pi e^{-\alpha}\left\{\frac{2}{r^{2}}\left(F_{24}\right)^{2}-e^{-\beta}\left(F_{14}\right)^{2}\right\}\right], \tag{3.7c}
\end{equation*}
$$

$$
\left(u^{4}\right)^{2}=\frac{(\varepsilon+p)^{-1}}{8 \pi} e^{-\alpha}\left[8 \pi p+e^{-\beta}\left(\frac{\beta_{1}}{r}-\frac{1}{r^{2}}\right)+\frac{1}{r^{2}}+8 \pi\left\{\frac{2}{r^{2}}\left(F_{24}\right)^{2}+e^{-\beta}\left(F_{14}\right)^{2}\right\}\right]
$$

$$
\left(F_{24}\right)^{2}=\left[e^{-\beta} \frac{\beta_{4}^{2}}{r^{2}}-e^{\alpha-2 \beta} \frac{a_{1} \beta_{1}}{r^{2}}-e^{\alpha}\left\{2(8 \pi) c^{-\alpha-\beta}\left(F_{14}\right)^{2}-e^{-\alpha}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)-\right.\right.
$$

$$
\left.-e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}^{2}}{4}-\frac{a_{11}}{2}+\frac{1}{r^{2}}\right)+\frac{1}{r^{2}}\right\}\left\{-\frac{1}{r^{2}}-2(8 \pi) e^{-\alpha-\beta}\left(F_{\lrcorner 4}\right)^{2}\right.
$$

$$
\left.\left.+e^{-a}\left(-\frac{\beta_{44}}{2}+\frac{\beta_{4}^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\alpha_{1}-\beta_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}^{2}}{4}-\frac{\alpha_{11}}{2}+\frac{1}{r^{2}}\right)\right\}\right]
$$

$$
\div \frac{2}{r^{2}}(8 \pi)\left[\begin{array}{c}
e^{-\beta}  \tag{3•8}\\
r \\
r
\end{array}\left(\alpha_{1}+\beta_{1}\right) \pm 2 e^{\frac{-\alpha-\beta}{2}} \frac{\beta_{4}}{r}\right]
$$

Thus the algebraic gravitational field equations enable us to express $\varepsilon, p, u^{1}, u^{4}$ and $F_{24}$ all in terms of $F_{14}, \alpha$ and $\beta$. Now $F_{14}$ will be determined from Maxwell's equations in terms of $\alpha, \beta$ and when suitable conditions of state are imposed on
$\varepsilon$ and $p, \alpha, \beta$ can be determined from $(3.7 a, b)$ $\varepsilon$ and $p, \alpha, \beta$ can be determined from ( $3 \cdot 7 a, b$ )

When in ( $3 \cdot 3 a, b$ ) and ( $3.5 a$ ) positive sign is taken, Maxwell's equations ( 2.8 ) reduce to

$$
e^{\frac{\beta-a}{2}}\left(F_{24,3}-F_{34,2}\right)+F_{23,1}=0
$$

$$
\begin{align*}
& \left(e^{\frac{\beta-\alpha}{2}} F_{34}\right)_{, 4}+F_{34,1}-F_{14,2} \sin \theta=0, \\
& \left(e^{\beta \cdot \underline{a}} F_{34}\right)_{4}+F_{34,1}-F_{14,3}=0, \\
& -F_{24,3}+F_{34,2}+F_{23,4}=0
\end{align*}
$$

From ( $3.9 d$ ) and $(3.9 c)$ we get

$$
F_{14,9}-F_{14,2} \sin \theta=0,
$$

which leads to the solution

$$
F_{14} \equiv F_{14}\lceil\{\phi-\log (\cot \theta+\operatorname{cosec} \theta)\}, r, t] .
$$

Similarly other solutions arising out of different combinations of signs in $(3.3 a, b)$ and ( $3.5 a$ ) can be obtained.
Case (ib) Distribution of Discrete Particles.
In the case of discrete particle distribution we have

$$
p=0
$$

and $(3.7 a),(3 \cdot 8)$ reduce to

$$
\begin{align*}
8 \pi e^{-\alpha-\beta}\left(F_{14}\right)^{2}= & e^{-a}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}{ }^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}{ }^{2}}{4}-\frac{\alpha_{11}}{2}\right), \\
\left(F_{24}\right)^{2}= & {\left[e^{-\beta^{\beta_{4}{ }^{2}}} \frac{r^{2}}{r^{2}}-e^{\alpha-2 \beta} \frac{\alpha_{1} \beta_{1}}{r^{2}}+e^{\alpha}\left\{8 \pi e^{-\alpha-\beta}\left(F_{14}\right)^{2}+\frac{1}{r^{2}}\left(1-e^{-\beta}\right)\right\}^{2}\right.} \\
& \left.+e^{\alpha}\left\{8 \pi e^{-\alpha-\beta}\left(F_{14}\right)^{2}+\frac{1}{r^{2}}\left(1-e^{-\beta}\right)\right\} \frac{\left.\beta_{1}-\alpha_{1}\right)}{r}\right] \\
& \div \frac{2}{r^{2}}(8 \pi)\left[e^{-\beta\left(\frac{\left.\alpha_{1}+\beta_{1}\right)}{r} \pm 2 e^{\frac{-\alpha-\beta}{2}} \frac{\beta_{4}}{r}\right] .}\right.
\end{align*}
$$

From (3.10) and (3.11) we conclude that $F_{i t}$ and $F_{24}$ are functions of $r$ and $t$. So Maxwell's equations will become

$$
\begin{align*}
& \left(r^{2} e^{\frac{-a-\beta}{2}} F_{14}\right)_{, 1} \sin \theta-e^{\frac{\beta-a}{2}} F_{24} \cos \theta=0, \\
& \pm\left(e^{\frac{\beta-a}{2}} F_{24}\right)_{, 4}+F_{24,1}=0, \\
& \left(r^{2} e^{\frac{-a-\beta}{2}} F_{14}\right)_{, 4} \sin \theta+F_{24} \cos \theta=0 .
\end{align*}
$$

Equations ( $3 \cdot 12 a$ ) and ( $3 \cdot 12 c$ ) will hold at all points only when
and

$$
\left(r^{2} e^{\frac{-a-\beta}{2}} F_{14}\right)_{, 1}=0,
$$

$$
F_{24}=0
$$

But we started with the assumption that none of the $F_{i j}{ }^{\text {'s }}$ is zero.
Hence discrete particle distribution does not permit all the $F_{i j}$ 's to be non-zero. Gase (iia)

$$
F_{14}=F_{29}=0
$$

Then

$$
\left.\begin{array}{l}
E_{1}^{1}=-E_{4}^{4}=-\frac{e^{-\alpha}}{r^{2}}\left[\left(F_{24}\right)^{2}+\left(F_{34}\right)^{2} \operatorname{cosec}^{2} \theta\right] \\
E_{2}^{2}=E_{3}^{s}=0, \\
E_{4}^{1}= \pm \frac{e^{\frac{-\alpha-\beta}{2}} r^{2}}{2}\left[\left(F_{24}\right)^{2}+\left(F_{3_{4}}\right)^{2} \operatorname{cosec}^{2} \theta\right] .
\end{array}\right\}
$$

Equations (3.1), $(3 \cdot 2 a),(2 \cdot 3 a, c, d)$ and (3.13) will give

$$
\begin{align*}
& -8 \pi p=e^{-a}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}{ }^{2}}{4}-\frac{a_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\beta_{1}-\alpha}{2 r}{ }^{1}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}^{2}}{4}-\frac{a_{11}}{2}\right), \quad(3 \cdot 14 a), \\
& 8 \pi \varepsilon=\frac{2}{r^{2}}\left(1-e^{-\beta}\right)-e^{-\alpha}\left(\frac{\beta_{44}}{4}+\frac{\beta_{1}^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{a_{1}^{2}-a_{1} \beta_{1}}{4}+\frac{\alpha_{11}}{2}\right),(3 \cdot 14 b) \\
& \left(u^{1}\right)^{2}=\frac{(\varepsilon+p)^{-1}}{8 \pi} e^{-\beta}\left[-8 \pi p-\frac{1}{r^{2}}+e^{-\beta}\left(\frac{1}{r^{2}}+\frac{\alpha^{1}}{r}\right)+\frac{8 \pi e^{-\alpha}}{r^{2}}\left\{\left(F_{24}\right)^{2}+\left(F_{34}\right)^{2} \operatorname{cosec}^{2} \theta\right\}\right](3 \cdot 14 c) \\
& \left(u^{4}\right)^{2}=\frac{(\varepsilon+p)^{-1}}{8 \pi} e^{-\alpha}\left[+8 \pi p+\frac{1}{r^{2}}+i^{\beta}\left(\frac{\beta_{1}}{r}-\frac{1}{r^{2}}\right)+\frac{8 \pi e^{-\alpha}}{r^{2}}\left\{\left(F_{24}\right)^{2}+\left(F_{34}\right)^{2} \operatorname{cosec}^{2} \theta\right\}\right] \\
& \frac{8 \pi}{r^{2}}\left\{\left(F_{24}\right)^{2}+\left(F_{34}\right)^{2} \operatorname{cosec}^{2} \theta\right\}=\left[e^{-\beta} \frac{\beta_{4}^{2}}{r^{2}}+e^{a}\left\{\frac{1}{r^{2}}+e^{-\beta}\left(\frac{\beta_{1}}{r}-\frac{1}{r^{2}}\right)+8 \pi p\right\}\left\{\frac{1}{r^{2}}\right.\right. \\
& \left.\left.-e^{-\beta}\left(\frac{1}{r^{2}}+\frac{\alpha_{1}}{r}\right)+8 \pi p\right\}\right] \div\left[\frac{e^{-\beta}}{r}\left(\alpha_{1}+\beta_{1}\right) \pm 2 e^{\frac{-\alpha-\beta}{2}} \frac{\beta_{4}}{r}\right] .
\end{align*}
$$

and equations ( $2 \cdot 8$ ) will imply

$$
\left.\begin{array}{l}
F_{24,3}-F_{34,2}=0 \\
F_{24,1} \pm\left(e^{\frac{\beta-a}{2}} F_{24}\right)_{, 4}=0 \\
F_{34,1} \pm\left(e^{\frac{\beta-a}{2}} F_{34}\right)_{, 4}=0
\end{array}\right\}
$$

Case (iib) For discrete particle distribution, i.e., $p=0(3 \cdot 14 a)$ is specified by

$$
\alpha=\beta=f(r+t)+g(r-t),
$$

and $(3 \cdot 14 b, c, d)$ and $(3 \cdot 15)$ will turn out to be

$$
\begin{gather*}
8 \pi \varepsilon=\frac{2}{r^{2}}\left(1-e^{-f-g}\right), \\
\left.\left(u^{1}\right)^{2}=\frac{e^{-f-g}}{8 \pi \varepsilon}\left[-\frac{1}{r^{2}}+e^{-f-g}\left(\frac{1}{r^{2}}+\frac{f^{\prime}+\dot{g}}{r}\right)+\frac{e^{f+g}}{4 \dot{g}}\left\{\frac{\left(1-e^{-f-g}\right)^{2}}{r^{3}}-\frac{4}{r} e^{-2(f+g}\right) f^{\prime} \dot{g}\right\}\right]
\end{gather*}
$$

$$
\begin{align*}
& \left(u^{4}\right)^{2}=\frac{e^{-f-g}}{8 \pi \varepsilon}\left[\frac{1}{r^{2}}+e^{-f-g}\left(\frac{f^{\prime}+\dot{g}}{r}-\frac{1}{r^{2}}\right)+\frac{e^{f+g}}{4} \dot{g}\left\{\frac{\left(1-e^{-f-g}\right)^{2}}{r^{3}}-\frac{4}{r} e^{-2(f+g)} f^{\prime} \dot{g}\right\}\right] \\
& \frac{8 \pi}{r^{2}}\left\{\left(F_{24}\right)^{2}+\left(F_{3_{4}}\right)^{2} \operatorname{cosec}^{2} \theta\right\} \equiv \frac{e^{2(f+g)}}{4 f^{\prime}(\text { or } 4 \dot{g})}\left[\frac{\left(1-e^{-f-g}\right)^{2}}{r^{3}}-\frac{4}{r} e^{-2(f+g)} f^{\prime} \dot{g}\right]
\end{align*}
$$

where

$$
f^{\prime} \equiv \frac{d \dot{f}}{d(r+t)}, \dot{g} \equiv \frac{d g}{\bar{d}(\underline{r-t})} .
$$

Maxwell's equations (3.16) lead to

$$
\begin{aligned}
& F_{24,3}-F_{34,2}=0, \\
& F_{24} \equiv F_{24}\{(r \pm t), \theta, \phi\}, \\
& F_{34} \equiv F_{34}\{(r \pm t), \theta, \phi\},
\end{aligned}
$$

Further if we take $F_{34}=0$ (3.16) will give

$$
F_{24} \equiv F_{24}\{(r \pm t), \theta\}
$$

If $F_{24} \equiv F_{24}\{(r+t), \theta\}$ then $\dot{g}=0$ and right hand side of $(3 \cdot 18 b, c)$ will become infinite. So we take

$$
F_{24} \equiv F_{24}\{(r-t), \theta\}
$$

which implies $f^{\prime}=0$, i.e., $f=$ constant which can be taken to be zero without loss of generality. And $(3 \cdot 18 a, b, c)$ and $(3 \cdot 19)$ are

$$
\begin{gathered}
8 \pi \varepsilon=\frac{2}{r^{2}}\left(1-e^{-g}\right), \\
\left(u^{3}\right)^{2}=\frac{e^{-g}}{\theta \pi \varepsilon}\left[-\frac{1}{r^{2}}+e^{-g}\left(\frac{1}{r^{2}}+\frac{\dot{g}}{r}\right)+\frac{e^{g}}{4 \dot{g}}\left(\frac{\left(1-e^{-g}\right)^{2}}{r^{3}}\right]\right. \\
\left(u^{4}\right)^{2}=\frac{e^{-g}}{8 \pi \varepsilon}\left[\frac{1}{r^{2}}+e^{-g}\left(\frac{\dot{g}}{r}-\frac{1}{r^{2}}\right)+\frac{e^{g}}{4 \dot{g}} \frac{\left(1-e^{-g}\right)^{2}}{r^{2}}\right], \\
8 \pi\left(F_{12}\right)^{2}=8 \pi\left(F_{241} 1^{2}=\frac{e^{2 g}}{4 \dot{g}} \frac{\left(1-e^{-g}\right)^{2}}{r} .\right.
\end{gathered}
$$

In this case the current vector $J^{i}$ calculated from (2.9) iurns out to be zero.
C.se (iiia) Spherically Symmetric Electromagnetic Distribution.

In case (i) and (ii) the electromagnetic field has not been assumed to be necessarily spherically symmetric although the material distribution is spherically symmetric. If we impose conditions of spherical symmetry on $F_{i j}$, from the considerations of groups of motions we find that only two of $F_{i j}{ }^{s}-F_{14}, F_{23}$ are none-zero. Then

$$
\begin{align*}
& F_{12}=F_{13}=0, \\
& E_{1}^{1}=-E_{2}^{2}=-E_{3}^{s}=E_{4}^{4}=\frac{1}{2}\left[e^{-a-\beta}\left(F_{14}\right)^{2}+\frac{\left(F_{13}\right)}{r^{4}} \operatorname{cosec}^{2} \theta\right]^{2}, \\
& E_{4}^{1}=0
\end{align*}
$$

Equations (3.1), (3.2a), (2.3a, c, d) and (3.20) give

$$
\begin{align*}
&-8 \pi p=- 8 \pi E_{1}^{1}+e^{-a}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}{ }^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\beta_{1} \alpha_{1}-\alpha_{1}{ }^{2}}{4}-\alpha_{11}\right)(3 \cdot 21 a) \\
& 8 \pi \varepsilon=3(8 \pi) E_{1}^{1}+\frac{2}{r^{2}}\left(1-e^{-\beta}\right)-e^{-\alpha}\left(\frac{\beta_{41}}{2}+\frac{\beta_{4}{ }^{2}}{4}-\frac{\alpha_{1} \beta_{4}}{4}\right) \\
&+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\alpha_{1}{ }^{2}-\alpha_{1} \beta_{1}}{4}+\frac{\alpha_{11}}{2}\right), \\
&\left(u^{1}\right)^{2}= \frac{\left(\varepsilon+p^{-1}\right.}{8 \pi} e^{-\beta}\left[-8 \pi p-8 \pi E_{1}^{1}-\frac{1}{r^{2}}+e^{-\beta}\left(\frac{1}{r^{2}}+\frac{\alpha_{1}}{r}\right)\right], \\
&\left(u^{4}\right)^{2}= \frac{(\varepsilon+p)^{-1}}{8 \pi}-e^{-a}\left[8 \pi p+8 \pi E_{\perp}^{1}+\frac{1}{r^{2}}+e^{-\beta}\left(\frac{\beta_{1}}{r}-\frac{1}{r^{2}}\right)\right. \\
& 4(8 \pi) E_{1}^{1}=-\left\{\frac{e^{-\beta}}{r}\left(\beta_{1}-\alpha_{1}\right)-\frac{2}{r^{2}}\left(e^{-\beta}-1\right)-2 A\right\} \pm\left[\left\{\frac{e^{-\beta}}{r}\left(\beta_{1}-\alpha_{1}\right)+\frac{2}{r^{2}}\left(1 \cdot e^{-\beta}\right)-2 A\right\}^{2}\right. \\
&-4\left\{e^{-\alpha-\beta} \frac{\beta_{4}^{2}}{r^{2}}+e^{\alpha-\beta}\left[\frac{e^{-\beta}}{r}\left(\alpha_{1}-\beta_{1}\right)\left(A+\frac{e^{-\beta}}{r^{2}}-\frac{1}{r^{2}}\right)+\frac{e^{-2 \beta}}{r^{4}}\right.\right. \\
&\left.\left.\left.-2 \frac{e^{-\beta}}{r^{4}}+\frac{1}{r^{4}}-e^{-2 \beta} \frac{\alpha_{1} \beta_{1}}{r^{2}}+A^{2}+2 \frac{A}{r^{2}}\left(e^{-\beta}-1\right)\right]\right\}\right]^{1 / 2},
\end{align*}
$$

where

$$
A=e^{-a}\left(\frac{\beta_{41}}{2}+\frac{\beta_{4}^{2}}{4}-\frac{a_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}^{2}}{4}-\frac{\alpha_{11}}{2}\right) .
$$

Equations (2:8) imply

$$
\left.\begin{array}{l}
F_{14} \equiv F_{14}(r, t) \\
F_{23} \equiv F_{23}(\theta, \phi)
\end{array}\right\}
$$

Denoting the right hand side of $(3 \cdot 22)$ by $B(r, t)$, we obtain in view of $(3 \cdot 20)$

$$
4(\varepsilon \pi)\left(F_{23}\right)^{2}=r^{4}\left[2 B-4(8 \pi) e^{\left.\left.-\alpha-\beta F_{14}\right)^{2}\right] \sin ^{2} \theta, ~}\right.
$$

which by virtue of (3.23) implies

$$
r^{4}\left[2 B-4(8 \pi) e^{-\alpha-\beta}\left(F_{14}\right)^{2}\right]=k_{1}(\text { constant }),
$$

or

$$
4(8 \pi)\left(F_{14}\right)^{2}=e^{\alpha+\beta}\left(\left(2 B-\frac{k_{1}}{r^{4}}\right)\right.
$$

In particular if we assume $(3 \cdot 17)$ to hold in this case, then $(3 \cdot 21 a, b, c, d)$ and ( $3 \cdot 22$ ) will become

$$
8 \pi p=8 \pi E_{1}^{1}=-\frac{1}{2 r^{2}}\left(1-e^{-f-g}\right) \pm \frac{e^{-f-g}}{r}\left(f^{\prime} \dot{g}\right)^{1 / 2}
$$

[ 380 ]

$$
\begin{align*}
& 8 \pi \varepsilon=3(8 \pi) p+\frac{2}{r^{2}}\left(1-e^{-f-g}\right) \\
& \left(u^{1}\right)^{2}=\frac{(\varepsilon+p)^{-1}}{8 \pi} e^{-f-g}\left[-\frac{2}{r} e^{-f-g}\left(f^{\prime} \dot{g}\right)^{1 / 2}+e^{-f-g} \frac{\left(f^{\prime}+\dot{g}\right)}{r}\right], \\
& \left(u^{4}\right)^{2}=\frac{(\varepsilon+p)^{-1}}{8 \pi}-e^{-f-g}\left[\frac{2}{r} e^{-f-g}\left(f^{\prime} \dot{g}\right)^{1 / 2}+e^{-f-g} \frac{\left(f^{\prime}+\dot{g}\right)}{r}\right],
\end{align*}
$$

which for $p>0, \varepsilon>0, \varepsilon>3 p,\left(u^{1}\right)^{2}>0$ and $\left(u^{4}\right)^{2}>0$ imply

$$
\begin{equation*}
e^{f+g}>0, e^{-f-g}\left(f^{\prime} \dot{g}\right)^{1 / 2}>\frac{1}{2 r}\left(1-e^{-f-g}\right), \frac{f^{\prime}+}{2}>\left(f^{\prime} \dot{g}\right)^{1 / 2} \tag{327}
\end{equation*}
$$

Case (iiii) In the case of discrete particle distribution (3.21a,b,c,d) and (3.22) lead to

$$
\begin{align*}
& 8 \pi E_{1}^{1}=e^{-\alpha}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{a_{1} \beta_{1}-a_{1}^{2}}{4}-\frac{\alpha_{11}}{2}\right)\right. \\
& 8 \pi \varepsilon=2(8 \pi) E_{1}^{1}+\frac{2}{r^{2}}\left(1-e^{-\beta}\right)-\frac{e^{-\beta}}{r}\left(a_{1}-\beta_{1}\right), \\
&\left(u^{1}\right)^{2}=\frac{e^{-\beta}}{8 \pi \varepsilon}\left[-8 \pi E_{1}^{1}-\frac{1}{r^{2}}+e^{-\beta}\left(\frac{1}{r^{2}}+\frac{\alpha_{1}}{r}\right)\right], \\
&\left(u^{4}\right)^{2}=\frac{e^{-\alpha}}{8 \pi \varepsilon}\left[8 \pi E_{1}^{1}+\frac{1}{r^{2}}+e^{-\beta}\left(\frac{\beta_{1}}{r}-\frac{1}{r^{2}}\right)\right], \\
& {\left[8 \pi E_{1}^{1}+\frac{1}{r^{2}}-e^{-\beta}\left(\frac{\alpha_{1}}{r}+\frac{1}{r^{2}}\right)\right]\left[8 \pi E_{1}^{1}+\frac{1}{r^{2}}+e^{-\beta}\left(\frac{\beta_{1}}{r}-\frac{1}{r^{2}}\right)=-e^{-\alpha-\beta} \frac{\beta_{4}^{2}}{r^{2}},\right.}
\end{align*}
$$

and

$$
8 \pi\left(F_{29}\right)^{2}=k_{2} \sin ^{2} \theta
$$

where
$k_{2}=r^{4}\left[2 e^{-\alpha}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}{ }^{2}}{4}-\frac{a_{4} \beta_{4}}{4}\right)+2 \epsilon^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}{ }^{2}}{4}-\frac{a_{11}}{2}\right)-8 \pi e^{-\alpha-\beta}\left(F_{14}\right)^{2}\right]$, $8 \pi\left(F_{14}\right)^{2}=e^{\alpha+\beta}\left[2 e^{-\alpha}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}{ }^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+2 e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{a_{1} \beta_{1}-a_{1}{ }^{2}}{4}-\frac{\alpha_{12}}{2}\right)-\frac{k_{2}}{r^{2}}\right]$.

Further in the absence of matter we have

$$
\varepsilon=p=0
$$

and equations ( $2 \cdot 3 a, b, c, d$ ) are

$$
\begin{align*}
-8 \pi E_{1}^{1}= & \frac{1}{r^{2}}-e^{-\beta}\left(\frac{1}{r^{2}}+\frac{a_{1}}{r}\right) \\
8 \pi E_{1}^{ \pm}= & e^{-\alpha}\left(\frac{\beta_{44}}{2}+\frac{\beta_{4}{ }^{2}}{4}-\frac{\alpha_{4} \beta_{4}}{4}\right)+e^{-\beta}\left(\frac{\beta_{1}-\alpha_{1}}{2 r}+\frac{\alpha_{1} \beta_{1}-\alpha_{1}{ }^{2}}{4}-\frac{\alpha_{11}}{2}\right] \\
-8 \pi E_{\mathrm{I}}^{\mathrm{T}}= & \frac{1}{r^{2}}+e^{-\beta}\left(\frac{\beta_{1}}{r}-\frac{1}{r^{2}}\right) \\
& \beta_{4}=0
\end{align*}
$$

Equations (3.32) will lead to the solution

$$
e^{\beta}=e^{-a}=\left(1-\frac{a}{r}-\frac{b}{r^{2}}\right)^{-1}
$$

and

$$
4 \pi\left[\left(F_{14}\right)^{2}+\frac{\left(F_{23}\right)^{2}}{r^{4}} \operatorname{cosec}^{2} \theta\right]=\frac{b}{r^{4}}
$$

If we take here $F_{23}=0$, Nordström solution is obtained as a particular case.

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# On the Rate of Decay of Solutions of Neutral Functional Differential Systems 

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1. In recent years differential systems involving rather arbitrary retarded arguments have been studied more and more intensively. This is due to the importance of such systems in various branches of technology. Recently, Driver [1, 2], Lakshmikantham and Lella [4] and others have dealt with the neutral functional differential systems, and have obtained several remarkable results. However, all the theory developed concerns itself with the stability and boundedness problems of neutral functional differential systems by considering one-sided estimates of solutions. In many cases of practical applications, it is natural to expect that, an estimation of the lower bound for the rate at which the solutions approach the origin or the invariant set, would yield useful refinements of stability notions. In the present paper upper and lower bounds are determined for the rate of growth or decay of solutions of the neutral func ional differential systems which arise in one approach in the twobody problem of classical electrodynamics. These bounds are obtained by using comparison principle and the extension of Lyapunov's method. We introduce the concepts of relative-strict-stability and relative-strict-boundedness of the neutral functional differential systems and the two sided estimates obtained ensure that the relative motion remains in tube-like domains
2. Consider the two systems of neutral functional differential equations

$$
\begin{align*}
& \quad x^{\prime}(t)=f\left(t, x(t), x(g(t, x(t))), x^{\prime}(g(t, x(t)))\right), t>t_{0} \\
& \text { with } x(t)=\phi(t) \text { on }\left[a, t_{0}\right], \\
& y^{\prime}(t)=h\left(t, y(t), y(g(t, y(t))), y^{\prime}(g(t, y(t)))\right), \quad t>t_{0} \\
& \text { with } y(t)=\psi(t) \text { on }\left[a, t_{0}\right],
\end{align*}
$$

where $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right) ; y=\left(y_{1}, y_{2}, \ldots, y_{n}\right) ; f=\left(f_{1}, f_{2}, \ldots, f_{n}\right)$; $h=\left(h_{1}, h_{2}, \ldots, h_{n}\right)$ and $g=\left(g_{1}, g_{2}, \ldots, g_{m}\right)$. For $\alpha<t_{0}, \alpha \leqslant g_{j}(t, x(t))<t$ and $a \leqslant g_{j}(t, y(t))<t, j=1,2, \ldots, m$, for all $t \geqslant t_{0}$ and $x^{\prime}\left(g_{j}(t, x(t))\right)$ and $\left.y^{\prime} g_{j}(t, y(t i))\right)$ stands for $\frac{d x}{d \epsilon}$ and $\frac{d y}{d t}$ evaluated at $g_{j}(t, x(t))$ and $g_{j}(t, y(t))$ respectively. The word 'Neutral' is used because the equations could be considered both for retarded
arguments as well as for advanced arguments. However, in these systems the retarded arguments will depend upon the dependent variables as well as the independent variable.

Suppose $f, h$ and $g$ are continuous functions of their arguments in their appropriate domains that is $f$ is continuous in $t, x(t), x(g(t, x(t))), x^{\prime}(g(t, x(t)))$ in the domain $D=D^{1+n+2 n m}$ (an open connected set) in $E^{1+n+2 n m}$, an ( $1+n+2 n m$ ) dimensional Euclidean space. A function $x(t)$ is said to be a solution of $(2 \cdot 1)$ if it satisfies the following conditions:
(i) $(t, x(t), x(g(t, x(t)))) \varepsilon D-G$
for each compact set $G \subset D, t_{0} \leqslant t<\beta$,
(ii) $x(t)=\phi(t), \alpha \leqslant t \leqslant t_{0}$,
(iii) $x^{\prime}(t)=f\left(t, x(t), x(g(t, x(t))), x^{\prime}(g(t, x(t)))\right)$

$$
\text { for almost all } t \varepsilon\left(t_{0}, \beta\right) \text {. }
$$

Similarly we can define $y(t)$ as the solution of (2.2).
The existence and uniqueness of Neutral Functional-Differential Equations have been discussed by many authors. Driver [1,2] has given references of many results obtained by various authors in this direction.

Let $x(t)$ and $y(t)$ be any two solutions of (2•1) and (2.2) respectively. Let $\sigma[x, y]=\|x-y\|$ denote the distance between two points $x$ and $y$. Our interest lies in the following definitions. The systems $(2 \cdot 1)$ and $(2 \cdot 2)$ are said to be
$\left(S_{1}\right)$ relatively-strictly-equi-stable, if for each $\varepsilon_{1}>0, t_{0} \geqslant 0$, it is possible to find positive functions $d_{1}=d_{1}\left(t_{0}, \varepsilon_{1}\right), d_{2}=d_{2}\left(t_{0}, \varepsilon_{1}\right)$ and $\varepsilon_{2}=\varepsilon_{2}\left(t_{0}, \varepsilon_{1}\right)$ which are continuous in $t_{0}$ for each $\varepsilon_{1}$, $\varepsilon_{2}<d_{2} \leqslant d_{1}<\varepsilon_{1}$ such that

$$
\varepsilon_{2}<\sigma[x(t), y(t)]<\varepsilon_{1}, t>t_{0},
$$

whenever

$$
\left.d_{2} \leqslant v[\phi(t),\}^{\prime}(t)\right] \leqslant d_{1}, \alpha \leqslant t \leqslant t_{0} ;
$$

$\left(S_{2}\right)$ relatively-strictly-equi-norm-bounded, if for each $\alpha_{1}>0, \mathrm{t}_{0} \geqslant 0$, it is possible to find for every $a_{2}$, satisfying $a<a_{2}<a_{1}$, positive functions $\beta_{1}=\beta_{1}\left(i_{u}, \alpha_{1}\right)$, $\beta_{2}=\beta_{2}\left(t_{0}, a_{1}, a_{2}\right)$, which are cuntinuous in $t_{0}$ for each $\alpha_{1}, \beta_{2}<\beta_{1}, \beta_{2}<\alpha_{2}$ such that

$$
\beta_{2}<\sigma[x(t), y(t)]<\beta_{1}, t>t_{0},
$$

whenever

$$
a_{2} \leqslant \sigma[\phi(t), \psi(t)] \leqslant a_{1}, \quad \alpha \leqslant t \leqslant t_{0} .
$$

3. Let $V(t, x, y) \geqslant 0$ be defined and continuous function of its argument in its appropriate domain, and satisfies a Lipschitz condition in $x$ and $y$ for each $t$. Define, for small $\lambda>0$
(3:1) $V^{*}(t, x, y)=\lim _{\lambda \rightarrow 0^{+}} \sup _{\lambda} \frac{1}{\lambda}\left[V\left(t+\lambda, x(t)+\lambda f\left(t, x(t), x(g(t, x(t))), x^{\prime}(g(t, x(t)))\right)\right.\right.$,

$$
\begin{gathered}
\left.y(t)+\lambda h\left(t, y(t), y(g(t, y(t))), y^{\prime}(g(t, y(t)))\right)\right) \\
-V(t, x(t), y(t))] .
\end{gathered}
$$

We formulate the following lemma which is useful for our subsequent discussion.
Lemsa. Let the function $V^{*}(t, x, y)$ of $(3 \cdot 1)$ satisfy the inequality

$$
W_{u}(t, V(t, x, y)) \leqslant V^{*}(t, x, y) \leqslant W_{1}(t, V(t, x, y))
$$

where $W_{1}(t, r)$ and $W_{2}(t, r)$ are continuous real valued funclions for $t \varepsilon I=[0, \infty), r \geqslant 0$. Let $r(t)$ and $\rho(t)$ be the maximal and minimal solutions of the scalar differential equations

$$
\begin{array}{r}
r^{\prime}=W_{1}(t, r), r\left(t_{0}\right)=r_{0}, \\
\rho^{\prime}=W_{2}(t, \rho), \rho\left(t_{0}\right)=\rho_{o},
\end{array}
$$

respectively for all $t \geqslant t_{0}$. If $x(t)$ and $y(t)$ be any two solutions of (2•1) and (2.2) with their initial functions $\phi(t)$ and $\psi(t)$ respectively on $\left[a, t_{0}\right]$ such that

$$
\rho_{0} \leqslant V\left(t_{0}, \phi(t), \psi(t)\right) \leqslant r_{0}, a \leqslant t \leqslant t_{0}
$$

then

$$
\rho(t) \leqslant V(t, x(t), y(t)) \leqslant r(t), t>t_{0} .
$$

Proof. We shall prove the right half of the inequality ( 3.5 ), similar reasoning can be used for the left half. Define $m(t)=V(t, x(t) y(t))$. then $m\left(t_{0}\right) \leqslant r_{0}$. Further for small $\lambda>0$,

$$
\begin{aligned}
m(t+\lambda)-m(t) \leqslant & C\left[\left\|x(t+\lambda)-x(t)-\lambda f\left(t, x(t), x(g(t, x(t))), x^{\prime}(g(t, x(t)))\right)\right\|\right. \\
& \left.+\left\|y(t+\lambda)-y(t)-\lambda h\left(t, y(t), y(g(t, y(t))), y^{\prime}(g(t y(t)))\right)\right\|\right] \\
& +V\left(t+\lambda, x(t)+\lambda f\left(t, x(t), x(g(t, x(t))), x^{\prime}(g(t, x(t)))\right),\right. \\
& \left.y(t)+\lambda h\left(t, y(t), y(g(t, y(t))), y^{\prime}(g(t, y(t)))\right)\right)-V(t, x(t), y(t))
\end{aligned}
$$

where $C>0$ is a Lipschitz constant. This together with ( $3 \cdot 1$ ) and ( $3 \cdot 2$ ) implies the inequality

$$
\lim _{\lambda \rightarrow 0^{+}} \frac{1}{\lambda}[m(t+\lambda)-m(t)] \leqslant W_{1}(t, m(t)) .
$$

The standard argument used in [2] can now be followed to establish the desired right-half of the result in ( $3 \cdot 5$ ).

Corresponding to the definition $\left(S_{1}\right)$, we can formulate the definition $\left(S_{1}{ }^{*}\right)$ wi:h respect to (3.3) and (3.4).
$\left(S_{1}{ }^{*}\right)$ For each $\eta_{1}>n_{,} t_{0} \geqslant 0$, there exist positive functions $\delta_{1}=\delta_{1}\left(t_{0}, \eta_{1}\right)$, $\delta_{2}=\delta_{2}\left(t_{0}, \eta_{1}\right)$ and $\eta_{2}=\eta_{2}\left(t_{0}, \eta_{1}\right)$ which are continuous in $t_{0}$ for each $\eta_{1}, \eta_{2}<\delta_{2} \leqslant \delta_{1}<\eta_{1}$ such that

$$
\eta_{2}<\rho(t) \leqslant r(t)<\eta_{1}, t>t_{0}
$$

whenever

$$
\delta_{2} \leqslant \rho_{0} \leqslant r_{0} \leqslant \delta_{1}
$$

Definition $\left(S_{2}{ }^{*}\right)$ may be formulated similarly.
4. We list below certain assumptions which will be used subsequently.

$$
b\left(\sigma\left[x, y^{\prime}\right]\right) \leqslant V(t, x, y) \leqslant a(\sigma[x, y])
$$

where $a(r), b(r)$ are continuous, strictly increasing functions for $r \geqslant 0$ with $a(0)=b(0)=0$,

$$
b(r) \rightarrow \infty \text { as } r \rightarrow \infty .
$$

The following theorems provide sufficient conditions for relative-strictstability and relative-strict-boundedness of the systems ( $2 \cdot 1$ ) and ( $2 \cdot 2$ ).
Theorem 1. Let the assumptions of the lemma hold, together with (4*1). Then ( $S_{1}{ }^{*}$ ) implies that the systems $(2 \cdot 1)$ and (2.2) are relatively-strictly-equi-stable.

Proof. Let $\varepsilon_{1}>0$ and $t_{0} \geqslant 0$ be given. Since ( $S_{1}{ }^{*}$ ) holds, given $\eta_{1}=b\left(\varepsilon_{1}\right)>0$, $t_{0} \geqslant 0$, there exist positive functions $\delta_{1}=\delta_{1}\left(t_{0}, \eta_{1}\right), \delta_{2}=\delta_{2}\left(t_{u}, \eta_{1}\right)$ and $\eta_{2}=\eta_{2}\left(t_{0}, \eta_{1}\right), \eta_{2}<\delta_{2} \leqslant \delta_{1}<\eta_{1}$,
such that

$$
\eta_{2}<\rho(t) \leqslant r(l) \leqslant \eta_{1}, t>t_{0}
$$

whenever

$$
\delta_{2} \leqslant \rho_{0} \leqslant r_{0} \leqslant \delta_{1} .
$$

Choose an $\varepsilon_{2}>0$ such that $a\left(\varepsilon_{2}\right) \leqslant \eta_{2}$ and $\varepsilon_{2}<\varepsilon_{1}$. Let $x(t)$ and $y(t)$ be any two solutions of $(2 \cdot 1)$ and $(2 \cdot 2)$ respectively, such that,

$$
\rho_{0} \leqslant V\left(t_{0}, \phi(t), \psi(t)\right) \leqslant r_{0}, \text { for } \alpha \leqslant t \leqslant t_{0}
$$

Then it follows from the lemma that

$$
\rho(t) \leqslant V(t, x(t), \gamma(t)) \leqslant r(t), \quad t>t_{0} .
$$

Further $(4 \cdot 1),(4 \cdot 5)$ and $(4 \cdot 4)$, show that there exist two positive functions $d_{1}=d_{1}\left(t_{0}, \varepsilon_{1}\right), d_{2}=d_{2}\left(t_{0}, \varepsilon_{1}\right)$ such that

$$
d_{2} \leqslant \sigma[\phi(t), \psi(t)] \leqslant d_{1} \text { implies } \delta_{2} \leqslant V\left(t_{0}, \phi(t), \psi(t)\right) \leqslant \delta_{1}
$$

for $a \leqslant t \leqslant t_{0}$, and vice versa. Thus whenever $d_{2} \leqslant r[\phi(l), \psi(t)] \leqslant d_{1}$, it follows from the assumptions of the lemma that ( $4 \cdot 6$ ) is true. Now we claim that $\varepsilon_{2}<\sigma[x(t), y(t)]<\varepsilon_{1}$, for $t>t_{0}$, provided $d_{2} \leqslant \sigma[\phi(t), \psi(t)] \leqslant d_{1}$, for $\alpha \leqslant t \leqslant t_{0}$.

Suppose on the contrary, there exist solutions $x(t)$ and $y(l)$ of (2.1) and (2.2) respectively, satisfying $d_{2} \leqslant \sigma[\phi(t), \psi(t)] \leqslant d_{1}$, such that for some $t=t_{1}>t_{0}$,
$\sigma\left[\lambda\left(t_{1}\right), \jmath\left(t_{1}\right)\right]=\varepsilon_{1}$ or $v\left[\lambda\left(t_{1}\right), y\left(t_{1}\right)\right]=\varepsilon_{2}$. In the first instance, using the inequalities in $(4 \cdot 1),(46)$ and (4.3) we arrive at the contradiction

$$
b\left(\varepsilon_{1}\right) \leqslant V\left(t_{1}, x\left(t_{1}\right), y\left(t_{1}\right)\right) \leqslant r\left(t_{1}\right)<b\left(\varepsilon_{1}\right) .
$$

On the other hand, if $\sigma\left[x\left(t_{1}\right), y\left(t_{1}\right)\right]=\varepsilon_{2}$, we arrive to a similar contradiction

$$
a\left(\varepsilon_{2}\right) \geqslant V\left(t_{1}, x\left(t_{1}\right), y\left(t_{1}\right)\right) \geqslant \rho\left(t_{1}\right)>\eta_{2} \geqslant a\left(r_{2}\right),
$$

because of the inequalities in $(4 \cdot 1),(4 \cdot 6)$ and $(4 \cdot 3)$. Which proves that the condition $\left(S_{1}{ }^{*}\right)$ implies ( $\bar{S}_{1}$ ).
Theorem 2. Let the assumptions of the lemma hold, together with (4•1) and (4•2). Then $\left(S_{2}{ }^{*}\right)$ implies that the systems $(2 \cdot 1)$ and $(2 \cdot 2)$ are relatively-strictly-equi-normbounded.
Proof. Let $\alpha_{1}>0$ and $t_{0} \geqslant 0$ be given. Choose $\alpha_{2}$ such that $0<\alpha_{2} \leqslant a_{1}$. Let $\phi(t)$ and $\psi(t)$ be such that $\alpha_{2} \leqslant \sigma[\phi(t), \psi(t)] \leqslant \alpha_{1}$, for $\alpha \leqslant t \leqslant t_{0}$. Then because of (41) it is possible to find positive numbers $\hat{a}_{1}=\hat{a}_{1}\left(\alpha_{1}\right)$ and $\hat{a}_{2}=\hat{a}_{2}\left(a_{2}\right)$ such that

$$
\hat{a}_{2} \leqslant V\left(t_{0}, \phi(t), \psi(t)\right) \leqslant \hat{a}_{1}, a \leqslant t \leqslant t_{0}
$$

Let $\left(S_{2}{ }^{*}\right)$ holds. Then given $\hat{a}_{1}>0, t_{0} \geqslant 0$, there exist, for every $\hat{a}_{2}, 0<\hat{a}_{2} \leqslant \hat{a}_{1}$, two positive functions $l_{1}=l_{1}\left(t_{0}, \hat{a}_{1}\right), l_{2}=l_{2}\left(t_{0}, \hat{a_{1}}, \hat{a_{2}}\right), l_{2}<l_{1}, l_{2}<\hat{a_{2}}$ such that

$$
l_{2}<\rho(t) \leqslant r(t)<l_{1}
$$

whenever

$$
\hat{a}_{2} \leqslant \rho_{0} \leqslant r_{0} \leqslant \hat{a}_{1} .
$$

Since $b(r) \rightarrow \infty$ as $r \rightarrow \infty$, it is possible to find positive functions $\beta_{1}=\beta_{1}\left(t_{0}, \alpha_{1}\right)$, $\beta_{2}=\beta_{2}\left(t_{0}, a_{1}, a_{2}\right)$, such that

$$
l_{1} \leqslant b\left(\beta_{1}\right), l_{2} \geqslant a\left(\beta_{2}\right), \beta_{2}<a_{2}, \beta_{2}<\beta_{1} .
$$

Now supposing the contrary and proceeding as in the proof of Theorem 1 , one can prove that $\left(S_{2}\right)$ is true.
5. Corresponding to the systems $(2 \cdot 1)$ and (2.2), we consider the systems

$$
\begin{align*}
x^{\prime}(t)= & f\left(t, x(t), x(g(t, x(t))), x^{\prime}(g(t, x(t)))\right) \\
& \quad+F\left(t, x(t), x(g(t, x(t))), x^{\prime}(g(t, x(t)))\right) \\
\text { with } x(t) & =\phi(t) \text { on }\left[\alpha, t_{0}\right] \\
y^{\prime}(t)= & h\left(t, y(t), y(g(t, y(t))), y^{\prime}(g(t, y(t)))\right) \\
& \quad+H\left(t, y(t) \cdot y(g(t, y(t))), y^{\prime}(g(t, y(t)))\right) \\
\text { with } y(t) & =\psi(t) \text { on }\left[\alpha, t_{0}\right]
\end{align*}
$$

respectively, where $F$ and $H$ are perturbed functions and obtain similar properties under constantly acting perturbations.

Assume that

$$
\begin{align*}
& \left\|F\left(t, x(t), x(g(t, x(t))), x^{\prime}(g(t, x(t)))\right)\right\| \\
& +\left\|H\left(t, y(t), y(g(t, y(t))), y^{\prime}(g(t, y(t)))\right)\right\| \\
& \leqslant n V(t, x, y), \quad(n>0) .
\end{align*}
$$

If the solutions of the systems $(5 \cdot 1)$ and (5.2) satisfy the definitions $\left(S_{1}\right)$ and $\left(S_{3}\right)$, whenever $(5 \cdot 3)$ is satisfied, we shall say that the systems ( $2 \cdot 1$ ) and (2.2) satisfy the definitions $\left(S_{1}\right)$ and $\left(S_{2}\right)$ weakly. Then it is easy to prove analogous results for weak-relative-strict-stability and weak-relative-strict-boundedness. We merely state the following two theorems.

Theorem 3. Let the assumptions of the lemma hold except that the condition (3.2) is replaced by

$$
W_{2}(t, V(t, x, y)) \div V^{*}(t, x, y)+\alpha V(t, x, y) \leqslant W_{1}(t, V(t, x, y))
$$

where $\alpha=n C$. Assume that $(4 \cdot 1)$ holds. If the condition $\left(S_{1}{ }^{*}\right)$ holds, then the systems (2•1) and (2.2) are weakly-relatively-strictly-equi-stable.

ThEOREM 4. Let the assumptions of the lemma hold except that the condition (3.2) is replaced by $(5 \cdot 4)$. Assume that $(4 \cdot 1)$ and $(4 \cdot 2)$ hold. If the condition $\left(S_{2}{ }^{*}\right)$ holds, then the systems $(2 \cdot 1)$ and $(2 \cdot 2)$ are weakly-relatively-strictly-equi-norm-bounded.

We note that many properties including the stability and boundedness studied in [4] satisfied by a scalar differential equation imply the corresponding properties satisfied by the neutral functional differential systems.

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## On order of partial sum of Laplace series

## By

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1. Let $f(\theta, \phi)$ be a function defined for the range $0 \leqslant \theta \leqslant \pi, 0 \leqslant \phi \leqslant 2 \pi$, and integrable in the sense of Labesgue on the surface of the sphere $S$. The Laplace series corresponding to this function is

$$
f(\theta, \phi) \sim \frac{1}{2 \pi} \Sigma\left(n+\frac{1}{2}\right) \int_{s} \int f\left(\theta^{\prime}, \phi^{\prime}\right) P_{n}(\cos \omega) d \sigma^{\prime}
$$

where

$$
\cos \omega=\cos \theta \cos \theta^{\prime}+\sin \theta \sin \theta^{\prime} \cos \left(\phi-\phi^{\prime}\right)
$$

and

$$
d \sigma^{\prime}=\sin \theta^{\prime} d \sigma^{\prime} d \phi^{\prime}
$$

and the Legendre polynomials $P_{n}(x)$ is defined by the relation

$$
\left(1-2 x t+t^{2}\right)^{-t}=\sum_{n=0}^{\infty} t^{n} P_{n}(x)
$$

We define on the lines of Kogbetliantz ${ }^{3}$, by $f(\omega)$ the function,

$$
f(\omega)=\frac{1}{2 \pi \sin \omega} \int_{c_{\omega}} f\left(\theta^{\prime}, \phi^{\prime}\right) d \sigma^{\prime}
$$

where the integral is taken along the small circle $c_{\omega}$ whose centre is the point $(\theta, \phi)$ on the surface of the sphere $S$ and whose curvilinear radius is $\omega$. The series ( $1 \cdot 1$ ) then reduces to the form

$$
\Sigma\left(n+\frac{1}{2}\right) \int_{0}^{\pi} f(\omega) \sin \omega P_{n}(\cos \omega) d \omega
$$

We write in this paper

$$
F(\omega)=f(\omega)(\sin \omega)^{\frac{1}{k}-k}
$$

and denote by $S_{n}$ the $n$th partial sum of the series $(1 \cdot 1)$. We also write

$$
\left\{\frac{d}{d x}\left[P_{n}(x)\right]+\frac{d}{d x}\left[P_{n+1}(x)\right]\right\}_{x=\cos \omega}=P_{n}^{\prime}(\cos \omega)+P_{n+1}^{\prime}(\cos \omega) .
$$

2. For Fourier series of a function $F(x)$ integrable in the interval $(-\pi, \pi)$ and periodic with period $2 \pi, S z a s z^{5}$ has established the following theorcm on the order of partial sums.

Theorem:

$$
\text { If } \int_{0}^{h}|F(x)| d x=0\left(\frac{h}{\log \frac{h}{t}}\right) \text {, }
$$

as $h \rightarrow 0$,
then

$$
\sum_{m=0}^{n} A_{m} \cos m x=0(\log \log n)
$$

where

$$
A_{n}=\frac{1}{\pi} \int_{0}^{2 \pi} F(x) \cos n x d x
$$

The order of partial sums of Legendre series has been calculated by Wilsone who proved the following
Theorem.

$$
(2 \cdot 2)
$$

$$
\text { For } 0 \leqslant k \leqslant \frac{1}{2} \text {, }
$$

if the integral

$$
\int_{0}^{\pi} f(\cos \theta)(\sin \theta)^{k+k} d \theta
$$

exists, then

$$
S_{n}(\cos \theta) \equiv \sum_{r=0}^{n} a_{\nu} P_{\nu}(\cos \theta)=\left[\begin{array}{l}
0(n),\left(0<k \leqslant \frac{1}{2}\right) \\
0(\log n),(k=0)
\end{array}\right.
$$

Later on in 1966 Beohar ${ }^{1}$ has proved on the order of partial sum of Legende series the following theorem:
(2.3) Theorem: At any interior point $x$ of the interval $(-1,1)$ for which

$$
\phi(t) \equiv \int_{0}^{t}|\phi(u)| d u=0\left\{\frac{t}{\left(\log \frac{t}{t}\right)^{\alpha}}\right\}
$$

as $t \rightarrow 0$,
then

$$
S_{n}(x)-f(x)=\left[\begin{array}{ll}
0(\log n)^{1-\alpha}, & (0 \leqslant \alpha<1) \\
0(\log \log n), & (\alpha=1)
\end{array}\right.
$$

Here we prove the following theorem on the order of partial sum of Laplace series.
(2•4) Theorem. If $\quad F(\omega) \in \operatorname{Lip}^{*}\left(\frac{1}{2}-k\right)$,
then

$$
S_{n}=\left[\begin{array}{ll}
0\left(n^{k}\right) & ,\left(0<k<\frac{1}{2}\right) \\
0(\log n), & (k=0)
\end{array}\right.
$$

3. In order to prove the theorem we need the following lemmas :

Lemma 1. ${ }^{2}$

$$
P_{n}^{\prime}(\cos \omega)+P_{n+1}^{\prime}(\cos \omega)=O\left(n^{2}\right) .
$$

for all values of $\omega$.

Lemma 2. ${ }^{2}$

$$
\left\{P_{n}{ }^{\prime}(\cos \omega)+P_{n+1}^{\prime}(\cos \omega)\right\} \sin \omega=O(n \sin \omega)
$$

in the range $\pi-n^{-1} \leqslant \omega \leqslant \pi$.

$$
\begin{align*}
& \text { Lemma 3. }{ }^{2} \\
& \qquad\left[P_{n}^{\prime}(\cos \omega)+P_{n+1}^{\prime}(\cos \omega)\right] \sin \omega=-\frac{2}{\pi} \\
& R\left[( n + 1 ) e ^ { \lfloor [ ( n + 1 ) \omega + \frac { \pi } { 4 } ] } ( \operatorname { c o t } \frac { \omega } { 2 } ) ! \left\{\left[\frac{1}{2}(n+1)^{-\frac{1}{2}}+O\left(n^{-3 / 2}(\sin \omega)^{-1}\right\}\right.\right.\right. \\
& \left.+O\left(n^{-\frac{1}{2}}(\sin \omega)^{-\frac{1}{2}}\right)+O\left(n^{-1} \omega^{-1}(\sin \omega)^{-\frac{1}{2}}\right)+O\left(n^{-3 / 2}(\sin \omega)^{-3 / 2}\right)\right]
\end{align*}
$$

in $\alpha_{n} \leqslant \omega \leqslant \pi-n^{-1}$.
4 Proof of the theorem. From the definition we have

$$
S_{n}=\sum_{k=0}^{n}\left(k+\frac{1}{2}\right) \int_{0}^{\pi} f(\omega) P_{k}(\cos \omega) \sin \omega d \omega .
$$

It is know that [4, pp. 178]

$$
\sum_{k=0}^{n}(2 k+1) P_{n}(x)=\frac{d}{d x}\left[P_{n}(x)\right]+\frac{d}{d x}\left[P_{n+1}(x)\right]
$$

Hence

$$
\begin{align*}
& \sum_{k=0}^{n}(2 k+1) P_{n}\left(\cos \omega=\left\{\frac{d}{d x}\left[P_{n}(x)\right]+\frac{d}{d x}\left[P_{n+1}(x)\right]\right\}_{x=\cos \omega}\right. \\
& \quad=P_{n}^{\prime}(\cos \omega)+P_{n+1}^{\prime}(\cos \omega) .
\end{align*}
$$

In view of $(1 \cdot 5),(4 \cdot 1)$ and $(4 \cdot 3) S_{n}$ is given by

$$
S_{n}=\frac{1}{2} \int_{0}^{\pi} F(\omega)\left[P_{n}^{\prime}(\cos \omega)+P_{n+1}^{\prime}(\cos \omega)\right](\sin \omega)^{k+\frac{1}{2}} d \omega,
$$

$$
\begin{aligned}
=\frac{1}{2}\left(\int_{0}^{\alpha_{n}}+\int_{\alpha_{n}}^{\pi-n^{-1}}\right. & \left.+\int_{\pi-n^{-1}}^{\pi}\right) F(\omega)\left[P_{n}^{\prime}(\cos \omega)+P_{n+1}^{\prime}(\cos \omega)\right](\sin \omega)^{k+\frac{1}{2}} d \omega \\
= & I_{1}+I_{2}+I_{3}
\end{aligned}
$$

where

$$
\begin{align*}
& \alpha_{n}=\left[\begin{array}{l}
n^{k / 2-1},\left(0<k<\frac{1}{\frac{1}{2}}\right) \\
(\log n)^{\frac{3}{2}} n^{-1},(k=0)
\end{array}\right. \\
& \text { First, } \quad I_{1}=O\left(n^{2}\right) \int_{0}^{a} n O\left(\omega^{\frac{1}{-k}}\right) \omega^{k+\frac{k}{k}} d \omega \text {, } \\
& =\left[\begin{array}{l}
0\left(n^{k}\right),\left(0<k<\frac{1}{2}\right) \\
0(\log n),(k=0)
\end{array}\right.
\end{align*}
$$

by the use of Lemma 1 and the fact that

$$
\begin{array}{r}
F(\omega)=0\left(\omega^{-k}\right) \\
{\left[\begin{array}{ll}
391
\end{array}\right]}
\end{array}
$$

(4.6) Next,

$$
I_{3}=O(n) \int_{k-n^{-1}}^{\pi}(\sin \omega)^{2 / t}+\frac{1}{2} d \omega
$$

$=O(1)$ by the use of Lemma 2 and the boundedness
of $F(\omega)$.
Finally, by the use of Lemma 3 we have
$(4 \cdot 7)$

$$
\begin{aligned}
& I_{2}=- \frac{1}{\pi} R\left[\int_{a_{n}}^{\pi-n^{-1}}(n+1) e^{i[(n+1) \omega+\pi / 4]} F(\omega)(\sin \omega)^{k-\frac{1}{2}}\left(\cot \frac{\omega}{2}\right)^{\frac{1}{2}}\right. \\
&\left.\left\{\Gamma_{2}^{2}(n+1)^{-1 / 2}+O\left(n^{-3 / 2}(\sin \omega)^{-1}\right)\right\}\right] d \omega \\
&+\left(n^{-1 / 2}\right) \int_{\alpha_{n}}^{\pi-n^{-1}} F(\omega)(\sin \omega)^{k-1} d \omega+O\left(n^{-1 / 2}\right) \int_{a_{n}}^{\pi-n^{-1}} F(\omega) \omega^{-1}(\sin \omega)^{k-1} d \omega \\
&+\left(n^{-3 / 2}\right) \int_{a_{n}}^{\pi-n^{-1}} F(\omega)(\sin \omega)^{k-2} d \omega \\
&=I_{2 \cdot 1}+ I_{2^{2} \cdot 2}+I_{2 \cdot 3}+I_{2 \cdot 4}+I_{2 \cdot 5} \cdot \\
& I_{3 \cdot 2}=O\left(n^{-1 / 2}\right) \int_{\alpha_{n}}^{\pi-n^{-1}} F(\omega)\left(\cot \frac{\omega}{2}\right)^{1 / 2}(\sin \omega)^{k-3 / 2} d \omega, \\
&= O\left(n^{-1 / 2}\right) \int_{a_{n}}^{\pi / 2} \omega^{\frac{1}{2}-\hbar} \cdot \omega^{-\frac{1}{2}} \omega^{k-3 / 2} d \omega+O\left(n^{-1 / 2}\right) \int_{\pi / 2}^{\pi-n^{-1}}(\sin \omega)^{k \cdot 2} d \omega, \\
&= O(1) .
\end{aligned}
$$

$(4 \cdot 7 \cdot 1)$
$(4 \cdot 7 \cdot 2)$

$$
\begin{aligned}
I_{2 \cdot 3} & =O\left(n^{-1 / 2}\right) \int_{a_{n}}^{\pi-n^{-1}} F(\omega)(\sin \omega)^{k-1} d \omega \\
& =O\left(n^{-1 / 2}\right) \cdot \int_{a_{n}}^{\pi / 2} \omega^{\frac{1}{2}-k} \omega^{k-1} d \omega+O\left(n^{-1}\right) \int_{\pi / 2}^{\pi-n^{-1}}(\sin \omega)^{/ 6-1} d \omega \\
& =O(1)
\end{aligned}
$$

$(4 \cdot 7 \cdot 3)$

$$
\begin{aligned}
I_{2 \cdot 4} & =O\left(n^{-\frac{1}{2}}\right) \int_{\alpha_{n}}^{\pi-n^{-1}} F(\omega) \omega^{-1}(\sin \omega)^{k-1} d \omega \\
& =O\left(n^{-\frac{1}{2}}\right) \int_{\alpha_{n}}^{\pi^{\prime 2}} \omega^{\frac{1}{2}-k} \omega^{-1} \omega^{k-1} d \omega+O\left(n^{-\frac{1}{2}}\right) \int_{\pi / 2}^{\pi-n^{-1}}(\sin \omega)^{k-1} d \omega, \\
& =O(1) .
\end{aligned}
$$

$(4 \cdot 7 \cdot 4)$

$$
\begin{aligned}
I_{2 \cdot 5} & =O\left(n^{-9 / 2}\right) \int_{\alpha_{n}}^{\pi-n^{-1}} F(\omega)(\sin \omega)^{k-2} d \omega \\
& =O\left(n^{-3 / 2}\right) \int_{\alpha_{n}}^{\pi / 2} \omega^{\frac{1}{2}-k_{\omega}} \omega^{k-2} d \omega+O\left(n^{-3 / 2}\right) \int_{\pi / 2}^{\pi-n^{-1}}(\sin \omega)^{k-2} d \omega \\
& =O(1)
\end{aligned}
$$

$$
I_{2 \cdot 1}=R\left[\begin{array}{c}
1 \\
\sqrt{\pi}
\end{array} e^{i \pi / 4}(n+1)^{\frac{y}{2}} \int_{\alpha_{n}}^{\pi n^{-1}} e^{i(n+1) \omega} F(\omega)\left(\cot _{\frac{1}{2}}^{m}\right)^{\frac{1}{2}}(\sin \omega)^{k-\frac{1}{k}} d \omega\right]
$$

Now

$$
\begin{gathered}
(n+1)^{\frac{1}{k}} \int_{a_{n}}^{\pi-n^{-1}} e^{i(n+1) \omega} F(\omega)\left(\cot \frac{\omega}{2}\right)^{\frac{1}{2}(\sin \omega)^{k-\frac{1}{2}} d \omega} \\
=\frac{子}{}(n+1)^{\frac{1}{2}}\left[\int_{a_{n}}^{\pi-n^{-1}} e^{i(n+1) \omega} F(\omega)\left(\cot \frac{\omega}{2}\right) \frac{1}{2}(\sin \omega)^{k-\frac{1}{2}} d \omega\right. \\
\left.=\int_{a_{n}-\mu_{n}}^{\pi-\mu_{n}-n^{-1}} e^{i(n+1) \omega} F\left(\omega+\mu_{n}\right)\left(\cot \frac{\omega+\mu_{n}}{2}\right)^{\frac{1}{2}}\left(\sin \left(\omega+\mu_{n}\right)\right)^{k-\frac{1}{2}} d \omega\right], \\
\left(\mu_{n}=\pi / n+1\right)
\end{gathered}
$$

and this is consequently less in modulus than

$$
\frac{1}{2}\left(J_{1}+J_{2}+J_{3}+J_{4}\right),
$$

$$
\begin{aligned}
& \text { where } \\
& J_{1}=(n+1)^{\frac{1}{4}} \int_{\alpha_{n}-\mu_{n}}^{\alpha_{n}} \left\lvert\, F\left(\omega+\mu_{n}\right)\left(\left.\cot \left(\frac{\omega+\mu_{n}}{2}\right)^{\frac{1}{2}}\left(\sin \left(\omega+\mu_{n}\right)\right)^{k-\frac{1}{2}} \right\rvert\, d \omega,\right.\right. \\
& J_{2}=(n+1)^{\frac{1}{2}} \int_{\pi-n^{-1}-\mu_{n}}^{\pi-n^{-1}}\left|F(\omega)\left(\cot \frac{\omega}{2}\right)^{\frac{1}{2}}(\sin \omega)^{k-\frac{1}{2}}\right| d \omega, \\
& J_{3}=(n+1)^{\frac{1}{2}} \int_{\alpha_{n}}^{\pi-n^{-1}-\mu_{n}}\left|F\left(\omega+\mu_{n}\right)-F(\omega)\right|\left(\cot \left(\frac{\omega+\mu_{n}}{2}\right)\right)^{\frac{1}{2}}\left(\operatorname{in}\left(\omega+\mu_{n}\right)\right)^{k-\frac{1}{2}} d \omega, \\
& J_{4}=(n+1)^{\frac{1}{2}} \int_{\alpha_{n}}^{\pi-n^{-1}-\mu_{n}}\left|\left(\cot \left(\frac{\omega+\mu_{n}}{2}\right)\right)^{\frac{1}{2}}\left(\sin \left(\omega+\mu_{n}\right)\right)^{k-\frac{1}{2}}-\left(\cot \frac{\omega}{2}\right)^{\frac{1}{2}}(\sin \omega)^{k-\frac{1}{2}}\right| \\
& F(\omega) d \omega,
\end{aligned}
$$

Now,

$$
\begin{align*}
J_{1} & =O\left(n^{\frac{1}{2}}\right) a_{n}^{1-k} \cdot a_{n}^{k-1} \cdot \mu_{n}  \tag{1}\\
& =O(1)
\end{align*}
$$

$\left(4 \cdot 7 \cdot 5_{2}\right) \quad J_{2}=O\left(n \frac{1}{}\right) \int_{\pi-n^{-1}-\mu_{n}}^{\pi-n^{-1}} \omega^{1}-k \omega^{k-1} d \omega$,

$$
=O(1)
$$

$\left(4 \cdot 7 \cdot 5_{8}\right) \quad J_{3}=O\left(n^{!}\right) \int_{\alpha_{n}}^{\pi}\left[0\left(\mu_{n}\right)^{k-k_{0}} \cdot \omega^{k \cdot 1}\right] d \omega$,

$$
=O\left(n^{l^{6}}\right)
$$

$\left(4 \cdot 7 \cdot 5_{4}\right) \quad J_{4}=O\left(n^{i}\right) \int_{a_{n}}^{\pi} \mu_{n} \omega^{k-2} \cdot \omega^{\frac{1}{2}-k} d \omega$, $=O(1)$.

Summing up we see that the theorem is proved for the case $0<k<\frac{1}{2}$.
To prove the theorem for the case $k=0$, it is sufficient to show that $J_{3}=0$ ( $\log n$ ), since the value of $k$ does not effect the other results.

Now when $k=0$, we have from $\left(4 \cdot 7 \cdot 5_{3}\right)$

$$
\begin{align*}
J_{3} & =0\left(n^{\frac{1}{2}}\right) \int_{\alpha_{r}}^{\pi} \mu_{n}^{4} \omega^{-1} d \omega \\
& =0(\log n)
\end{align*}
$$

Thus the theorem is completely proved

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# Effect of Viscosity on the Stability of Superposed Fluids 

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## Abstract

A study of the effect of viscosity on the stability of the superposed fluids is made and a dispersion relation is derived in case of two uniform fluids of constant densities and velocities. Kelvin-Helmoltz and Rayleigh-Taylor instabilities are discussed and it is found that these instabilities are independent of the viscosity.

## Introduction

S. Chandrasekher ${ }^{1}$ and recently Sharma and Srivastava ${ }^{2}$ have discussed the Kelvin-Helmoltz instability in the presence of a magnetic field for non viscous fluids. We derive here a dispersion relation for the superposed viscous fluids and discuss Kelvin-Helmoltz and Rayleigh-Taylor instabilities. It is first supposed that the fluids are in relative motion to obtain dispersion relation in case of Kelvin-Helmoltz instability and the fluids are set at rest to discuss RayleighTaylor instability.

The fluids are supposed to be viscous incompressible and moving with velocity $U$ in the $x$-direction and separated by a horizontal boundary $z=0$. Let $\delta \rho$ and $\delta_{p}$ be a small disturbance in density and pressure respectively such that the density at any point $(x, y, z)$ be $\rho+\delta \rho$ and pressure $p+\delta p$. The components of velocity in the perturbed state be $U+u, v, w$. The perturbed equation of fluids are :

$$
\begin{align*}
& \rho \frac{\partial u}{\partial t}+\rho U \frac{\partial u}{\partial x}+\rho w \frac{d U}{d z}=-\frac{\partial}{\partial x}(\partial p)+\mu \nabla u  \tag{l}\\
& \rho \frac{\partial v}{\partial t}+\rho U \frac{\partial v}{\partial x}=-\frac{\partial}{\partial y}\left(\delta_{p}\right)+\mu \nabla^{2} v \\
& \rho \frac{\partial w}{\partial t}+\rho U \frac{\partial w}{\partial x}=-\frac{\partial}{\partial z}(\delta p)-g^{\delta} \rho^{\prime}+\Sigma_{s} T_{s}\left[\left(\frac{\partial^{2}}{i x^{2}}+\frac{\partial^{2}}{\partial y^{2}}\right) \delta z_{s}\right] \delta\left(z-z_{s}\right) \\
& \quad+\mu \nabla^{2} w
\end{align*}
$$

where the coefficient of viscosity $\mu$ is supposed to be constant and the other symbols have their usual meaning.

Equation of continuity, for this case is

$$
\begin{equation*}
\frac{\partial u}{\partial x}+\frac{\partial v}{\partial y}+\frac{\partial w}{\partial z}=0 \tag{4}
\end{equation*}
$$

Also the density of every particle remains unchanged and $\delta z$ can be expressed in terms of the normal component of the velocity $w_{s}$ at $z_{s}$, corresponding perturba. tion equations are

$$
\begin{align*}
& \frac{\partial(\delta \rho)}{\partial t}+U \frac{\partial}{\partial x}(\delta \rho)=-w \frac{d \rho}{d z}  \tag{5}\\
& \frac{\partial\left(\delta z_{s}\right)}{\partial t}+U_{s} \frac{\partial}{\partial x^{-}}\left(\delta z_{s}\right)=w\left(z_{s}\right) \tag{6}
\end{align*}
$$

Derivation of dispersion relations
Following S. Chandrasekher we analyse the disturbance into normal modes dependent on $x, y$, and $t$, which is given by $\exp i\left(k_{x} x+k_{y} y+n t\right)$. For solutions having this dependence on $x, y$, and $t$, cquations (1) to (6) become,

$$
\begin{equation*}
i \rho n^{\prime} u+p(D U) w=-i k_{x} \delta p+\mu\left(D^{2}-k^{2}\right) u \tag{7}
\end{equation*}
$$

$$
i_{\rho n^{\prime} v}=-i k_{y} \delta p+\mu\left(D^{\prime \prime}-k^{\prime}\right) v
$$

$$
\begin{align*}
& i \rho n^{\prime} w=-D(\delta \rho)-g^{\delta \rho-k^{2}} \searrow I_{s} z_{s} \delta\left(z-z_{s}\right)+\mu\left(D^{2}-k^{3}\right) w  \tag{8}\\
& i k_{x} u+i k_{y} v+D w=0  \tag{9}\\
& i n^{\prime} \delta \rho=-w D \rho  \tag{10}\\
& i n^{\prime} \delta z_{s}=w\left(z_{s}\right) \tag{11}
\end{align*}
$$

where

$$
\begin{equation*}
D=\frac{d}{d z} ; n=n+k_{x} U \tag{12}
\end{equation*}
$$

Multiplying equation (7) by $-i k_{n}$ and equation (8) by $-i k_{y}$ and adding and then eliminating $\delta p$ between this equation and equation (9), with the help of equations (10) (11) and (12), we obtain

$$
\begin{align*}
& D\left[\rho n^{\prime} D w-\left\ulcorner k_{x}(D U) w 1+D\left\ulcorner i \mu\left(D^{2}-k^{2}\right) D u\right]-k^{2} \rho n^{\prime} w\right.\right.  \tag{13}\\
& \quad=g k^{\prime \prime}(D P) \frac{w}{n^{\prime}}+i k^{2} \mu\left(D^{2}-k^{2}\right) w-k^{4}\left[T_{s}\left(\frac{w}{n^{\prime}}\right)_{s} \delta\left(z-z_{s}\right)\right]
\end{align*}
$$

Here $w / n^{\prime}$ is continuous, hence integrating equation (13) between $z-\varepsilon$ and $z+\varepsilon$ and in the limit $\varepsilon \rightarrow 0$, we get

$$
\begin{equation*}
\Delta s\left[\rho n^{\prime} D w-\rho k_{x}(D U) w+i \mu\left(D^{2}-k^{\prime 2}\right) D w\right]=g k^{\prime 2}\left[\Delta_{s} \rho-\frac{k^{2}}{g} T_{s}\right]\left(\frac{\omega}{n^{\prime}}\right)_{s} \tag{14}
\end{equation*}
$$ where $\Delta_{s}$ has its usual meaning.

In case of constant density $\rho$ and constant streaming $U$, equation (13) becomes

$$
D\left[\rho n^{\prime} D w+i \mu\left(D^{2}-k^{2}\right) D w\right]-k^{2} \rho n^{5} w=i k^{2} \mu\left(D^{2}-k^{2}\right) w
$$

which can be written as

$$
\begin{equation*}
\left(D^{2}-k^{2}\right)\left(D^{2}-q^{2}\right) w=0 \tag{15}
\end{equation*}
$$

where

$$
q^{2}=k^{2}+\frac{i n^{\prime}}{\nu} ; v=\frac{\mu}{\rho} \text { (kincmatic viscosity). }
$$

The two fluids of densities $\rho_{1}$ and $\rho_{2}$ are separated by a horizontal boundary at $z=0$, and the density of the upper fluids is $\rho_{2}$ and since $w / n^{\prime}$ is continuous, $w$ cannot grow exponentially and further if the real part of $q$ is positive, the solutions for the two regions of the fluid, can be written as

$$
w_{1}=\left(A e^{l l^{l} z}+B e^{q_{1}} z\right) n_{1}^{\prime} \quad \text { for } z<0
$$

$$
w_{2}=\left(A e^{-k z}+B e^{-q_{2} z}\right) n_{2}^{\prime} \text { for } z>0
$$

where $n_{j}{ }^{\prime}=n+k_{x} U_{j} ; q_{j}^{2}=k^{2}+\frac{i n_{j}^{\prime}}{\nu_{j}}$ for $j=1,2$
$\mu\left(D^{2}+k^{2}\right) \frac{w}{n^{\prime}}$ is also continuous ${ }^{3}$, hence we can have

$$
\mu_{1}\left[2 k^{2} A+\left(q_{1}^{2}+k^{2}\right) B\right]=\mu_{\mathrm{zL}}\left[2 k^{2} A+\left(q_{2}^{2}+k^{2}\right) B\right]
$$

We obtain $B=\mathcal{N}$. $A$, where $\mathcal{N}=\frac{2 k^{2}\left(\mu_{1}-\mu_{2}\right)}{2 k^{2}\left(\mu_{2}-\mu_{1}\right)+i\left(n_{2}{ }^{\prime} \rho_{2}-n_{1}^{\prime}, \rho_{1}\right)}$
Thus the solutions of $w$ are given by

$$
\begin{aligned}
& w_{1}=A n_{1}{ }^{\prime}\left(e^{k z}+N e^{q_{1} z}\right) \quad \text { for } z<0 \\
& w_{2}=A n_{2}{ }^{\prime}\left(e^{-k z}+N e^{-q_{2} z}\right) \text { for } z>0
\end{aligned}
$$

Substituting $w_{1}$ and $w_{2}$ in equation (14), we obtain

$$
\rho_{2} n_{2}^{\prime 2}+\rho_{1} n_{1}^{\prime 2}=g k^{2}\left[\rho_{1}-\rho_{2}+\frac{k^{2}}{g} T\right][1+N]
$$

Putting the value of $\mathcal{N}$ from above we obtain, the dispersion relation

$$
\begin{align*}
& \left\{2 k^{2}\left(\mu_{2}-\mu_{1}\right)+i\left(\rho_{2} n_{2}^{\prime}-\rho_{1} n_{1}^{\prime}\right)\right\}\left(\rho_{21} n_{2}{ }^{\prime 2}+\rho_{1} n_{1}{ }^{\prime 2}\right)  \tag{16}\\
& =i g k^{2}\left[\rho_{1}-\rho_{2}+\frac{k^{2}}{g} T\right]\left[\rho_{2} n_{2}^{\prime}-\rho_{1} n_{1}^{\prime}\right]
\end{align*}
$$

## Discussion

We now discuss Kelvin-Helmolts and Rayleigh-Taylor instabilities.
(i) Kelvin-Helmoltz instability: In this we consider the cases of two liquids of same viscosity and of different viscosity.
(a) Two liquids of same viscosity: Here we have $\mu_{1}=\mu_{2}$, putting this condition in equation (16), we obtain

$$
\rho_{2} n_{2}^{\prime 2}+\rho_{1} n_{1}^{\prime 2}=g k^{2}\left[\rho_{1}-\rho_{2}+\frac{k^{2}}{g} T\right]
$$

For $n$ real we obtain the following stability condition

$$
\begin{equation*}
k_{. c}{ }^{2}\left(U_{1}-U_{2}\right)^{2}<\frac{g k}{a_{1} \alpha_{2}}\left\{\left(\alpha_{1}-\alpha_{2}\right)+\frac{k^{2} T}{g\left(\rho_{1}+\rho_{2}\right)}\right\} \tag{17}
\end{equation*}
$$

where

$$
a_{1}=\frac{\rho_{1}}{\rho_{1}+\rho_{2}} \text { and } \alpha_{2}=\frac{\rho_{2}}{\rho_{1}+\rho_{2}}
$$

(b) Two fluids of different viscosity: Here we have $\mu_{1} \neq \mu_{2}$, separating equation
(16) into real and imaginary parts, we obtain
or

$$
\begin{align*}
& 2 k^{2}\left(\mu_{2}-\mu_{1}\right)\left(\rho_{2} n_{2}^{\prime 2}+\rho_{1} n_{1}^{\prime 2}\right)=0  \tag{18}\\
& \rho_{2} n_{2}^{\prime 2}+P_{1} n_{1}^{\prime 2}=0
\end{align*}
$$

$$
\begin{equation*}
\rho_{2} n_{2}^{\prime 2}+\rho_{1} n_{1}^{\prime 2}=g k^{2}\left[\rho_{1}-\rho_{2}+\frac{k^{2}}{g} T\right] \tag{19}
\end{equation*}
$$

For $n$ to be real equation (18) gives $U_{1}=U_{2}$ and now putting this in equation (19) we obtain
either

$$
\rho_{1}>\rho_{2} \quad \text { or } \quad \alpha_{1}+\frac{k^{3} T}{g\left(\rho_{1}+\rho_{2}\right)}>a^{2}
$$

(20) i.e $\quad \rho_{1}>\rho_{2}$ or $k<k_{0}$
where $\quad k_{c}=\left\{\left(\rho_{2}-\rho_{1}\right) \frac{g}{T}\right\}^{1 / 3}$ and $k_{c}$ is called the critical value of $k$.
If we solve equation (19) for re.l value of $n$, we obtain the same condition as given in equation (17), putting this condition in (18) we obtain the stability condition for $U_{1} \neq U_{2}$

$$
\begin{equation*}
\rho_{2}>\rho_{1}+\frac{k^{2} T}{g} \text { or } k<k_{c} \tag{21}
\end{equation*}
$$

(ii) Rayleigh-T aylor instability: We put $U=0$ and then $n^{\prime}=n$ in equation (16) and since $w / n$ and $\mu\left(D^{2}+k^{2}\right) \frac{w}{n}$ are continuous, we obtain the dispersion relation for the case as

$$
\begin{equation*}
\left\{2 k^{2}\left(\mu_{12}-\mu_{1}\right)+i\left(\rho_{2}-\rho_{1}\right) n\right\}\left(\rho_{2}+\rho_{1}\right) n=i g k^{2}\left[\rho_{1}-\rho_{2}+\frac{k^{2}}{g}\right]\left(\rho_{2}-\rho_{1}\right) \tag{22}
\end{equation*}
$$

(a) Two liquids of same viscosity: Here we have $\mu_{1}=\mu_{2}$, equation (22) becomes

$$
n^{2}=g k^{2}\left[\alpha_{1}-a_{2}+\frac{k^{2} I}{\left(\rho_{1}+\rho_{2}\right) g}\right]
$$

For stability, $n$ must be real, which leads to following conditions
(i) either $\rho_{1}>\rho_{2}$ or (ii) $k<k_{c}$; where $k_{c}=\left\{\left(\rho_{2}-\rho_{1}\right) \frac{g}{T}\right\}^{1 / 2}$
which are the same as given by Chandrasekher ${ }^{4}$.
(b) Two liquids of different viscosity: Here we have $\mu_{1} \neq \mu_{2}$. Equating real and imaginary parts of the equation (22), we obtain stability condstion as given in case (ii), (a).

We therefore conclude that the stability of superposed fluids is independent of viscosity.

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## Series Involving Legendre Polynomials

## By

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1. Rainville [7] expressed Legendre polynomials of one variable in the form of Legendre polynomials of another variable as follows,

$$
P_{n}(\cos a)=\left(\frac{\sin \alpha}{\sin \beta}\right)^{n} \sum_{k=0}^{n}\binom{n}{k}\left[\frac{\sin (\beta-a)}{\sin \alpha}\right]^{n-k} P_{k}(\cos \beta) .
$$

Bhonsle [1] made the substitution $\beta=2 \alpha$ in the above relation and used it for evaluating a large number of interesting integrals.

An attempt has been made to obtain a similar series which has been used later on for evaluating certain integrals which are believed to be new. The following relation of Erdélyi [5, p. 264], has been considered for obtaining the series,

$$
\left(1-2 t x+t^{2}\right)^{-\frac{1+\nu}{2}} P_{\nu}^{\mu}\left[\frac{x-t}{\left(1-2 x t+t^{2}\right)^{\frac{1}{2}}}\right]=\sum_{r=0}^{\infty}\binom{r-\mu+r}{r} P_{\nu, r}^{\mu}(x) t^{r} .
$$

At the end, in the article 4 , an application to a heat conduction prcblem is considered.
2. In ( $1 \cdot 2$ ), substituting

$$
z=\cos \beta, t=\frac{\cos (\beta-a)}{\cos \alpha}
$$

and therefore,

$$
x-t=-\frac{\sin \beta \sin \alpha}{\cos \alpha} \text { and }\left(1-2 t x+t^{2}\right)=\frac{\sin ^{2} \beta}{\cos ^{2} \alpha},
$$

we obtain,

$$
\left[\frac{\sin \beta}{\cos \alpha}\right]^{-(1+\nu)} P_{v}^{\mu}(-\sin \alpha)=\sum_{r=0}^{\infty}\binom{v-\mu+r}{r} P_{\nu+r}^{\mu}(\cos \beta)\left[\frac{\cos (\beta-a)}{\cos a}\right]^{r}
$$

Let $\mu=m, v=m+n$, wet get,

$$
P_{m+n}^{m}(-\sin a)=\left[\frac{\sin \beta}{\cos \alpha}\right]^{m+n+1} \sum_{r=0}^{\infty}\binom{n+r}{r}\left[\frac{\cos (\beta-\alpha)}{\cos \alpha}\right]^{r} P_{m+n+r}^{m}(\cos \beta) .
$$

Let $m=0$ in (2:2), so that
$(2 \cdot 3) \quad P_{n}(-\sin \alpha)=\left(\frac{\sin \beta}{\cos \alpha}\right)^{n+1}{\underset{\Sigma}{n} 0}_{\infty}^{\infty}\binom{n+r}{r}\left[\frac{\cos (\beta-\alpha)}{\cos a}\right]^{r} P_{n+r}(\cos \beta)$.
But $P_{n}(-\sin \alpha)=(-1)^{n} P_{n}(\sin \alpha)$, therefore $(2 \cdot 3)$ becomes

$$
P_{n}(\sin \alpha)=(-1)^{n}\binom{\sin \beta}{\cos \alpha \alpha}^{n+1} \sum_{r=0}^{\infty}\binom{n+r}{r}\left[\frac{\cos (\beta-\alpha)}{\cos \alpha}\right]^{r} P_{n+r}(\cos \beta)
$$

Let $\beta=2 \alpha$ in $(2 \cdot 4)$ and we get,

$$
P_{n}(\sin \alpha)=(-1)^{n}(2 \cdot \sin \alpha)^{n+1}{\underset{y}{y}}_{\infty}^{r=0}\binom{n+r}{r} P_{n+r}(\cos 2 \alpha) .
$$

Writing $\cos 2 \alpha=x$ and $\sin a=\sqrt{\frac{-x}{2}}$, we have

$$
\left.\frac{(-1)^{n}}{2^{n+1}\left(\frac{1-x}{2}\right)^{\frac{n+1}{2}}} P_{n}\left(\sqrt{\frac{1-x}{2}}\right) \right\rvert\,=\sum_{r=0}^{\infty}\binom{n+r}{r} P_{n+r}(x)
$$

3. The following integrals have been evaluated, by making use of (2.6).

Using the orthogonal property $[4$, p. $170 ;(8),(10)]$ for Legendre polynomials, we have

$$
\text { (3.1) } \left.\quad \int_{-1}^{1} 2^{-\frac{n+1}{2}}(1-x)^{-\frac{n+1}{2}} P_{n}(\sqrt{1-x} 2)^{2}\right) P_{r}(x) d x\left\{\begin{array}{l}
=0 \quad, \quad r<n \\
=(-1)^{n}\binom{r}{r-n} \frac{2}{(2 r+1)}, r \geqslant n
\end{array}\right.
$$

If $p$ and $q$ are $+v e$ integers and $q \leqslant p$, we have the Adam's expansion [10, p. 331].

$$
P_{p}(\mu) P_{q}(\mu)=\sum_{r=0}^{q} \frac{A_{q-} A_{r} A_{p-r}}{A_{p+q-r}}\left(\frac{2 p+2 q-4 r+1}{2 p+2 q-2 r+1}\right) P_{p+q-2 r}(\mu)
$$

where $A_{r}=\frac{1.3 .5 \ldots \ldots(2 r-1)}{r!}$.
Thus making use of the result (3.2) with $(2 \cdot 6)$, we obtain, when $p+q-2 r \geqslant n$,

$$
\begin{gather*}
\int_{-1}^{1} \frac{(-1)^{n} P_{n}\left(\sqrt{\frac{1-x}{2}}\right)}{\left(\sqrt{\frac{1-x}{2}}\right)^{n+1}} P_{p}(x) P_{q}(x) d x=\sum_{r=0}^{q} \frac{A_{q-r} A_{r} A_{p+r \cdot} 2^{n+2}}{A_{p+q-r}(2 p+2 q-2 r+1)} \\
\binom{p+q-2 r}{p+q-2 r-n}=2^{2 n+1}\left(\frac{1}{2}\right)_{p}\left(\frac{1}{2}\right)_{q+q+1}(p+q)!(p+q)! \\
{[400]}
\end{gather*}
$$

$$
\times_{8} F_{5}\left[\begin{array}{c}
-p,-q,-\left(p+q+\frac{1}{2}\right),-\frac{(p+q-n)}{2},-\frac{1}{2}(p+q-n-1), \frac{1}{2} ; \\
\frac{1}{2}-p, \frac{1}{2}-q,-(p+q),-\frac{1}{2}(p+q),-\frac{1}{2}(p+q-1) ; \\
1
\end{array}\right] .
$$

Following the method of Shabde [8], if

$$
p_{1}+p_{2}+p_{3}+\cdots+p_{r}=s
$$

and

$$
\sqrt{\frac{1-x}{2}}=y, \text { we have, }
$$

$$
\begin{gather*}
\int_{-1}^{1} \frac{(-1)^{s}}{(2 y)^{s+r-1}} \cdot p_{p_{1}}(y) p_{p_{2}}(y) \ldots p_{p_{r}}(y) \cdot P_{s}\left(1-2 y^{2}\right) d y \\
\quad=\frac{(s!)^{3}}{(2 s+1)!} \prod_{t=1}^{r} \frac{\left(2 p_{t}\right)!}{\left(p_{t}\right)!\left(p_{t}\right)!}
\end{gather*}
$$

Grosswald's formula [1] for the $r$ th differentiation of $P_{n}(x)$ when $x=1$ is given by

$$
\left[\left(\frac{d}{d x}\right)^{r} P_{n}(x)\right]_{x=1}=\frac{(n+r)!}{2^{r} r!(n-r)!} .
$$

Thus differentiating (2•6) $r$ times with respect to $x$ and using the above formula, we have

$$
\begin{gather*}
\left.\left(\frac{d}{d x}\right)^{r}\left[\frac{(-1)^{n}}{\left(\sqrt{1-x} 2^{2-x}\right.}\right)^{n+1} P_{n}\left(\sqrt{\frac{1-x}{2}}\right)\right]_{x=1} \\
=2^{n+1-r}\binom{n+r}{r} \frac{(n+2 r)!}{n!r!}{ }_{3} F_{2}\left[\begin{array}{c}
1, n+r+1, n+2 r+1 ; 1 \\
r+1, n+1 ;
\end{array}\right]
\end{gather*}
$$

By using Neumann's formula [10, p. 320]

$$
Q_{n}(\mu)=\frac{1}{y} \int_{-1}^{1} \frac{P_{n}(y)}{\mu-y} d y
$$

with $(2 \cdot 6)$, we get

$$
\begin{equation*}
\int_{-1}^{1} \frac{(-1)^{n}}{\left(\frac{1-x}{2}\right)^{\frac{n+1}{2}}} \cdot \frac{P_{n}\left(\sqrt{\frac{1-x}{2}}\right)}{y-x} d x=2^{n+2} \sum_{k=0}^{\infty}\binom{n+k}{k} Q_{n+k}(x] . \tag{3.8}
\end{equation*}
$$

From Erdelyi [4, p. 171], when $r>-1$, we have

$$
\int_{0}^{1} x^{\sigma} P_{n}(x) d x=\frac{\left.\sqrt{ }^{-\pi} 2^{-\sigma-1} \Gamma_{1}^{\prime} 1+\sigma\right)}{\Gamma\left(1+\frac{\sigma}{2}-\frac{n}{2}\right) \Gamma\left(\frac{3}{2}+\frac{\sigma}{2}+\frac{n}{2}\right)} .
$$

Using the above result ( $3 \cdot 9$ ) with ( $2 \cdot 6$ ), we get

$$
\begin{gather*}
\int_{0}^{1} \frac{x^{\sigma}(-1)^{n}}{2^{n+1}\left(\frac{1-x}{2}\right)^{n+1}} P_{n}\left(\sqrt{\frac{1-x}{2}}\right){ }_{p} F_{q}\left[\begin{array}{l}
a_{1}, \alpha_{2}, \ldots, \alpha_{p} ; \pm x^{2} t^{2} \\
\beta_{1}, \beta_{2}, \ldots, \beta_{q} ;
\end{array}\right] d x \\
=\sum_{k=0}^{\infty}\binom{n+k}{k} \frac{\Gamma\left(\frac{1+\sigma}{2}\right) \Gamma\left(\frac{2+\sigma}{2}\right)}{2 \Gamma\left(1+\frac{\sigma}{2}-\frac{n+k}{2}\right)} \Gamma\left(\frac{3+\sigma+n+k}{2}\right) \\
{ }_{p+2} F_{q+2}\left[\begin{array}{l}
a_{1}, \alpha_{2}, \ldots \ldots, a_{p}, 1+\frac{\sigma}{2}, \frac{1}{2}+\frac{\sigma}{2} ; \\
\beta_{1}, \beta_{2}, \ldots \ldots, \beta_{q}, 1+\frac{\sigma}{2}-\frac{n+k}{2}, \frac{3+\sigma+n+k}{2} \pm t^{2} \\
\quad|t|<\mid \text { and } p=q+1
\end{array}\right.
\end{gather*}
$$

The hypergeometric series should be terminating when $p>q+1$. We can obtain many particular cases of the results (3.10).

We have due to Brafman [7, p. 168], the following generating function for Legendra polynomials,

$$
\begin{align*}
{ }_{2} F_{1}\left[\begin{array}{r}
a, 1-a \\
1 ;
\end{array} \begin{array}{rl} 
& \left.\frac{1-t-\rho}{2}\right] \times{ }_{2} F_{1}\left[\begin{array}{r}
a, 1-a ; 1+t-\rho \\
1 ;
\end{array}\right] \\
= & \sum_{n=0}^{\infty} \frac{(a)_{n}(1-a)_{n} t^{n}}{n!n!} P_{n}(x), \\
& \rho=\left(1-2 x t+t^{2}\right)^{\frac{1}{2}} .
\end{array} .\right.
\end{align*}
$$

where
Using the result ( $3 \cdot 11$ ) with $(2 \cdot 6)$, we get

$$
\begin{align*}
& \int_{-1}^{1} \frac{(-1)^{n}}{2^{n+1}\left(\frac{1-x}{2}\right)^{\frac{n+1}{2}}} P_{n}\left(\sqrt{\frac{1-x}{2}}\right)_{2} F_{1}\left[\begin{array}{r}
a, 1-a ; \frac{1-t-\rho}{2}
\end{array}\right] \times{ }_{2} F_{1}\left[\begin{array}{c}
a, 1-a ; 1+t \rho \\
1 ; 2
\end{array}\right] d x \\
& =\frac{(a)_{n}(1-a)_{n} \cdot 2 \cdot t^{n}}{n!n!(2 n+1)}{ }_{3} F_{2}\left[\begin{array}{cc}
a+n, 1-a+n, \frac{1}{2}+n ; \\
1+n, & \frac{8}{2}+n ;
\end{array}\right] .
\end{align*}
$$

If $a=-m$, we have
(3.13) $\int_{-1}^{1} \frac{(-1)^{n}}{2^{n+1}\left(\frac{i 1-x}{2}\right)^{\frac{n+1}{2}}} P_{n}\left(\sqrt{\frac{1-x}{2}}\right) P_{m}(t-\rho) P_{m}(t+\rho) d x$

$$
=\frac{(-m)_{n}(1+m)_{n} \cdot 2 \cdot t^{n}}{n!n!(2 n+1)} \cdot{ }_{3} F_{2}\left[\begin{array}{r}
-m+n, 1+m+n, n+\frac{1}{2} ; \\
n+1, n+\frac{3}{2} ;
\end{array}\right]
$$

We have from [7, p. 164(2)]

Using (3.14) with (2.6), we get

$$
\begin{gather*}
\int_{-1}^{1} \frac{(-1)^{n}(1-x t)^{n+b}}{2^{n+1}\left(\frac{1-x}{2}\right)^{\frac{1}{2}}} P_{n}\left(\sqrt{\frac{1-x}{2}}\right)_{2} F_{1}\left[\begin{array}{r}
\frac{b}{2}, \\
\frac{b+1}{2} ; t^{2}\left(x^{2}-1\right) \\
1
\end{array}\right] d x \\
\left.=\frac{(1-x t)^{2}}{n}\right] d\left(2 \left(\frac{b)_{n} t^{n}}{n+1)} \cdot{ }_{2} F_{1} \cdot\left[\begin{array}{r}
b+n, n+\frac{1}{2} ; \\
n+\frac{4}{2} ;
\end{array}\right]\right.\right.
\end{gather*}
$$

We have also due to Mitra [1],

$$
\int_{0}^{1} P_{n}\left(1-2 y^{2}\right) J_{0}(y x) y d y=\frac{1}{x} J_{2 n+1}(x) .
$$

Substituting $y$ for $\sqrt{\frac{1-x}{2}}$ in (2.6), we get

$$
\underset{2^{n+1} y^{n+1}}{(-1)^{n}} P_{n}(y)=\sum_{k=0}^{\infty}\binom{n+k}{k} P_{n+k}\left(1-2 y^{2}\right)
$$

Combining the above relation with (3.16), we get

$$
\begin{align*}
& \int_{0}^{1} \frac{(-1)^{n}}{2^{n+1} y^{n+1}} \cdot P_{n}(y) J_{0}(y x) y d y \\
& =\sum_{k=0}^{\infty}\binom{n+k}{k} \frac{1}{x} J_{2 n+2 k+1}(x)
\end{align*}
$$

Expanding $J_{o}(y x)$ in terms of $y x$ and integrating the left hand side of (3.17) term by term, with the help of (3.9), we get

$$
\begin{gather*}
\int_{0}^{1} \frac{(-1)^{n}}{2^{n+1} y^{n}} P_{n}(y) J_{0}(y x) d y \\
=\frac{(-1)^{n}}{2} \sum_{k=0}^{\infty} \frac{\left(\frac{1}{2}-\frac{n}{2}\right)_{k}\left(1-\frac{n}{2}\right)_{k}}{k!(1)_{k}(1-n)_{k}\left(\frac{3}{2}\right)_{k}}\left(-\frac{x^{2}}{4}\right)^{k}
\end{gather*}
$$

$$
=\frac{(-1)^{n}}{2}{ }_{2} F_{8}\left[\begin{array}{ll}
\frac{1}{2}-\frac{n}{2}, 1-\frac{n}{2} & : \\
1,1-n, \frac{x^{2}}{2} & ;
\end{array}\right]
$$

Thus equating (3.17) and (3.18), we get

$$
\frac{2}{x} \sum_{k=0}^{\infty}\binom{n+k}{k} J_{2 n+2 k+1}(x)=(-1)^{n}{ }_{2} F_{3}\left[\begin{array}{c}
\frac{1}{2}-\frac{n}{2}, 1-\frac{n}{2} ;-\frac{x^{2}}{4} \\
1,1-n, \frac{8}{2}
\end{array}\right]
$$

Using the following orthogonal property of Bessel functions, for $\nu>-1$,

$$
\int_{0}^{\infty} t^{-1} J_{v+2 n+1}(l) \cdot J_{v+2 m+1}(t) d t\left\{\begin{array}{l}
=0, m \neq n  \tag{3.20}\\
=\frac{1}{2(2 n+v+1)}, m=n
\end{array}\right.
$$

with (3.19), we get,

$$
\begin{gather*}
\int_{0}^{\infty}(-1)^{n}{ }_{2} F_{3}\left[\begin{array}{cc}
\frac{1}{2}-\frac{n}{2}, 1-\frac{n}{2} ; & -\frac{x^{2}}{4} \\
1,1-n, & \frac{3}{2} ;
\end{array}\right] J_{2 n+2 k+1}(x) d x \\
=\binom{n+k}{k} \frac{1}{(2 n+2 k+1)}
\end{gather*}
$$

Also we have,

$$
\begin{align*}
& \int_{0}^{\infty} x(-1)^{2 n}\left\{{ }_{2} F_{3}\left[\begin{array}{l}
\frac{1}{2}-\frac{n}{2}, 1-\frac{n}{2} ;-\frac{x^{2}}{4} \\
1,1-n, \frac{3}{2} ;
\end{array}\right]\right\}^{2} d x \\
& =\int_{0}^{\infty} \frac{4}{x}\left[\sum_{k=0}^{\infty}\binom{n+k}{k} J_{2 n+2 k+1}(x)\right]^{2} d x \\
& =2 \sum_{k=0}^{\infty} \frac{(n+k)!(n+k)!}{k!}, \frac{1}{(2 n+2 k+1)} \\
& =\frac{2}{(2 n+1)}{ }_{3} F_{2}\left[\begin{array}{r}
n+1, n+1, n+\frac{7}{3} ; \\
1, n+\frac{3}{2}:
\end{array}\right]
\end{align*}
$$

## 4. An application to a Heat Conduction problem.

Churchill [2] considered the problem given as follows. If $U(x, t)$ represents temperature in a non-homogeneous insulated bar, along the axis of $x$, with its ends at $x=-1$ and $x=+1$, in which the thermal conductivity is proportional to ( $1-x^{2}$ ), the heat equation takes the form,

$$
\frac{\partial U}{\partial t}=b \frac{\partial}{\partial x}\left[\left(1-x^{2}\right) \frac{\partial U}{\partial x}\right]
$$

where $b>0$ and constant, if the thermal coefficient $C_{0}$ is constant. The ends $x= \pm l$ are insulated because the conductivity vanishes there.

If $\quad U(x, 0)=f(x), \quad-1<x<1$
the solution of the above equation ( $4 \cdot 1$ ) is given by

$$
U(x, t)=\sum_{n=0}^{\infty} \frac{(2 n+1)}{2} \exp [-n(n+1) b t] P_{n}(x) \int_{-1}^{1} f(\xi) P_{n}(\xi) d \xi
$$

The equation $(4 \cdot 1)$ can also be solved with the application of Legendre transform. This transform was first introduced by Tranter ${ }^{[1]}$ in 1950 and was later on used by Churchill [3], and Normington and Blackwell [6] in solving the problem on heat transfer.

Let there be a function $f(x)$ which can be expanded in Legendre polynomials as follows,

$$
f(x)=\sum_{n=0}^{\infty} a_{n} P_{n}(x) .
$$

Legendre transforms are defined as

$$
\bar{f}(n)=\int_{-1}^{1} f(x) P_{n}(x) d x
$$

with the inversion theorem being given by

$$
f(x)=\frac{1}{3} \sum_{n=0}^{\infty}(2 n+1) \bar{f}(n) P_{n}(x) .
$$

Applying Legendre transform to (4•1) and making use of Legendre's equation,

$$
\frac{d}{d x}\left\{\left(1-x^{2}\right) \frac{d}{d x} P_{n}(x)\right\}=-n(n+1) P_{n}(x)
$$

we obtain finally,

$$
\frac{d}{d t} \bar{U}(n)=-b n(n+1) \bar{U}(n)
$$

The solution of $(4 \cdot 7)$ is given by

$$
\bar{U}(n)=A e^{-b n(n+1) t}
$$

The constant $A$ can be evaluated from the initial condition,

$$
U(x, 0)=f(x)
$$

Hence, $[\bar{U}(n)]_{t=0}=\bar{J}(n)$.
Therefore ( $4 \cdot 8$ ) becomes

$$
\bar{U}(n)=\bar{f}(n) e^{-b n(n+1) t}
$$

The inversion theorem (4.5) will therefore give

$$
U(x, t)=\frac{t}{2} \sum_{n=0}^{\infty}(2 n+1) \bar{f}(n) e^{-b n(n+1) t} P_{n}(x) .
$$

A function $f(x)$, which will suitably serve as initial temperature distributions is (2.6), expressed as

$$
f_{m}(x)=(-1)^{m} 2^{-\frac{m+1}{2}}(1-x)^{-\frac{m+1}{2}} P_{m}\left(\sqrt{\frac{1-x}{2}}\right) .
$$

From (3.1) we know that
$\int_{-1}^{1}(-1)^{m} 2^{-\frac{m+1}{2}}(1-x)-\frac{m+1}{2} P_{m}\left(\sqrt{\frac{1-x}{2}}\right) P_{n}(x) d x\left\{\begin{array}{l}=0, n<m \\ =\binom{n}{n-m} \frac{2}{(2 n+1)}, n \geq m .\end{array}\right.$
Therefore, we have

$$
\bar{f}_{m}(n)=\binom{n}{n-m} \frac{2}{(2 n+1)}, \text { for } n \geqslant m
$$

Thus from (4•10), the solution for the function $f(x)$ given by $(4 \cdot 11)$, will be,

$$
U_{m}\left(x_{1} t\right)=\sum_{n=0}^{\infty}\binom{n}{n-m} e^{-b n(n+1) t} P_{n}(x)
$$

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## A note on certain bilinear and bilateral generating relations for the Laguerre polynomial*

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## Abstract

The author presents here a systematic discussion on certain bilinear and bilateral generating relations associated with the generalized Leguerre polynomials $\left\{L_{n}^{(a)}(x)\right\}$ defined by

$$
(1-t)^{-1 \cdots \alpha} \exp \left(-\frac{x t}{1-t}\right)=\sum_{n=0}^{\infty} L_{i 2}{ }^{(\alpha)}(x) t^{n}
$$

It is observed that the four formulas proved recently by Jain [7] are only specialized or limiting forms of the results of Erdélyi [4], Meixner [8], and the author [12, 13, 14 ].

For the Laguerre polynomial

$$
L_{n}^{(a)}(z)=\sum_{k=0}^{n}\binom{n+a}{n-k} \frac{(-x)^{k_{-}}}{k!}
$$

Jain[7] has recently proved, in these Proceedings, one bilinear and three bilateral generating functions in the forms $\dagger$

$$
\begin{gather*}
\sum_{n=0}^{\infty} \frac{n!(1+b+k)_{n}}{(1+a)_{n}(1+b)_{n}} L_{n}^{(a)}(x) L_{n}^{(a)}(y) z^{n}  \tag{1}\\
=(1-z)^{-1-b-k} \exp \left[-\frac{(x+y) z}{1-z} \sum_{n=0}^{\infty} \frac{(1+b+k)_{n}}{n!(1+a)_{n}(1+b)_{n}}\left[\frac{x y z}{(1-z)^{2}}\right]^{n}\right. \\
\cdot{ }_{1} F_{1}\left[\begin{array}{l}
a-b-k ; \\
1+a+n ; \frac{x z}{1-z}
\end{array}\right]{ }_{1} F_{1}\left[\begin{array}{c}
-k ; y z \\
1+b+n ;
\end{array}\right] \\
{\underset{n=0}{\infty}{ }_{2} F_{1}\left[\begin{array}{l}
-n, d ; \\
1+a ;
\end{array}\right] L_{n}^{(b)}(y) z^{n}=(1-z)^{-1-b} \exp \left(-\frac{y z}{1-z}\right)}^{1+a} . \tag{2}
\end{gather*}
$$

[^0]\[

$$
\begin{align*}
& \cdot{ }_{3}{ }_{9}^{(1)}\left[\begin{array}{c}
\left.d, d, d, c, a-b ; c, 1+a, 1+a ;-\frac{x z}{1-z}, \frac{x z}{1-z}, \frac{x y z}{(1-z)^{2}}\right], \\
\sum_{n=0}^{\infty}{ }_{2} F_{1}\left[\begin{array}{c}
-n, d ; \\
1+a ;
\end{array}\right] L_{n}^{(b-n)}(y) z^{n}=(1+z)^{b} \exp (-y z) \\
\quad \cdot \Phi_{1}\left[d,-b ; 1+a ; \frac{x z}{1+z}, x y z\right]
\end{array}, \$\right. \text {, }
\end{align*}
$$
\]

and

$$
\sum_{n=0}^{\infty} \frac{(-)^{n}(-b)_{n}}{(1+a)_{n}} L_{n}^{(a)}(x)_{2} F_{1}\left[\begin{array}{c}
-n, d ;  \tag{4}\\
1+b-n ;
\end{array}\right] z^{n}
$$

$$
=(1+z)^{b}{ }_{3} \Phi_{M}^{(2)}\left[g, d,-b, d ; g, 1+a, 1+a ;-y z, \frac{x z}{1+z}, x y z\right],
$$

where $\Phi_{1}$ denotes one of Humbert's confluent hypergeometric functions in two arguments defined by means of (see, e.g., [5], p. 225)

$$
\Phi_{1}[a, \beta ; \gamma ; x, y]=\sum_{m, n=0}^{\infty} \underline{(\alpha)^{m+n}} \underset{(\gamma)_{m+n}}{(\beta)} \frac{x^{m}}{m!} \frac{\gamma^{n}}{n!},
$$

and

$$
\begin{aligned}
& 3^{\Phi_{G}^{(1)}}\left[\alpha, a, a, \beta, \beta^{\prime} ; \gamma, \gamma^{\prime}, \gamma^{\prime} ; x, y, z\right] \\
& =\sum_{m, n, p=0}^{\infty} \frac{(\alpha)_{m+n+p}+(\beta)_{m}\left(\beta^{\prime}\right)_{n}}{(\gamma)_{m}\left(\gamma^{\prime}\right)_{n+p}} \frac{x^{m}}{m!} \frac{y^{n}}{n!} \frac{z_{p}}{p!}, \\
& 3^{\Phi} \Phi_{M}^{(2)}\left[\alpha, \beta, \beta^{\prime}, \beta ; \gamma, \gamma^{\prime}, \gamma^{\prime} ; x, y, z\right] \\
& =\sum_{m, n, p=0}^{\infty} \frac{(a)_{m}(\beta)_{m+p}\left(\beta^{\prime}\right)_{n}}{(\gamma)_{m}\left(\gamma^{\prime}\right)_{r+p}} \frac{x^{m}}{m!} \frac{\gamma^{n}}{n!} \frac{z^{p}}{p!} .
\end{aligned}
$$

A simple analysis will reveal the fact that the confluent triple hypergeometric functions occurring on the right-hand sides of (2) and (4) are essentially the Humbert functions $\Phi_{1}$ and $\phi_{2}$, where $\Phi_{2}$ is defined by [5, p. 225]

$$
\Phi_{2}[\alpha, \beta ; \gamma ; x, y]=\sum_{m, n=0}^{\infty} \frac{(\gamma)_{m}(\beta)_{m}}{(\gamma)_{m+n}} \frac{x^{m}}{m!} \frac{y^{\prime \prime}}{n!},
$$

and consequently, the reduced forms of the formulas (2) and (4) are

$$
\begin{gather*}
\sum_{n=0}^{\infty}{ }_{22} F_{1}\left[\begin{array}{c}
-n, d ; \\
1+a ;
\end{array}\right] L_{n}^{(b)}(y) z^{n}  \tag{5}\\
=(1-z)^{d-b-1}(1-z+x z)^{-d} \exp \left(-\frac{y z}{1-z}\right) \\
\cdot \\
\Phi_{1}\left[d, a-b ; 1+a ; \frac{x z}{1-z+x z}, \frac{x y z}{(1-z)(1-z+x z)}\right]
\end{gather*}
$$

and

$$
\begin{align*}
& {\underset{S}{3}=0}_{\infty} \frac{(-)^{n} \cdot(-b)_{n}}{(1+a)_{n}} L_{n}^{(b)}(x)_{2} F_{1}\left[\begin{array}{cc}
-n, d & ; \\
1+b-n ;
\end{array}\right] z^{n}  \tag{6}\\
& =(1+z)^{b}(1+y z)^{-d} \Phi_{2}\left[-b, d ; 1+a ; \begin{array}{c}
x z \\
1+z, \frac{x y z}{1+y z}
\end{array}\right]
\end{align*}
$$

respectively.
The main purpose of the present note is to observe that the formulas (1) and (2) in much more general forms were proved many decades earlier by Erdélyi [4] and Meixner[8] and that all these resuls admit themselves of further elegant generalisation which we have derived elsewhere (see [12], [13] and [14]).

$$
L_{n}^{(\alpha)}(z)=\binom{n+\alpha}{n}{ }_{1} F_{1}\left[\begin{array}{c}
-n ; \\
1+\alpha ;
\end{array}\right],
$$

and the well-known Kummer transformation [5, p. 253]

$$
{ }_{1} F_{1}\left[\begin{array}{c}
a ; \\
c ;
\end{array}\right]=\exp (z){ }_{1} F_{1}\left[\begin{array}{c}
c-a ; \\
c ;
\end{array}\right]
$$

the formula (1) assumes the desired form

$$
\begin{align*}
& \underset{n=0}{\infty} \frac{(b+k)_{n}}{n!}{ }_{1} F_{1}\left[\begin{array}{c}
-n ; \\
a ;
\end{array}\right]{ }_{1} F_{1}\left[\begin{array}{c}
-n ; y \\
b ;
\end{array}\right] z^{n}=(1-z)^{-b-k} \sum_{n=0}^{\infty} \frac{(b+k)_{n}}{(a)_{n}(b)_{n} n!}  \tag{7}\\
& \cdot{ }_{1} F_{1}\left[\begin{array}{c}
b+k+n ; \frac{x z}{} \\
a+n ;{ }^{z}-1
\end{array}\right]{ }_{1} F_{1}\left[\begin{array}{c}
\left.b+k+n ; \frac{y z}{z+n}\right]\left[\begin{array}{c}
x y z \\
(1-z)^{2}
\end{array}\right]^{n} .
\end{array} .\right.
\end{align*}
$$

Obviously, this is a special case of Erdélyi's formula [4, p. 344]
(8)

$$
\begin{gathered}
\sum_{n=0}^{\infty} \frac{(\lambda)_{n}}{n!}{ }_{1} F_{1}\left[\begin{array}{c}
\rho-n ; \\
a ;
\end{array}\right]{ }_{1} F_{1}\left[\begin{array}{c}
a-n ; \\
b ; y
\end{array}\right] z^{n} \\
=(1-z)^{-\lambda} \sum_{n=0}^{\infty} \frac{(\lambda)_{n}}{(a)_{n}(b)_{n} n!} \Phi_{2}\left[\rho, \lambda+n ; a+n ; x, \frac{x z}{z-1}\right] \\
{[4091}
\end{gathered}
$$

$$
\Phi_{2}\left[r, \lambda+n ; b+n ; y, \frac{y z}{z-1}\right]\left[\frac{x y z}{(1-z)^{2}}\right]^{n},|z|<1
$$

when $n=\sigma=0$ and $\lambda=b+k$, Erdélyi [4, p. 344] does state the specia! case $\rho=\sigma=0$ of his formula (8) in the form*

$$
\begin{align*}
& \sum_{n=0}^{\infty} \frac{(\lambda)_{n} n!}{\left.(a)_{n}(b)\right)_{n}} L_{n}^{(a-1)}(x) L_{n}^{(b-1)}(y) z^{n}=(1-z)^{-\lambda} \cdot \sum_{n=0}^{\infty} \frac{(\lambda)_{n}}{(a)_{n}(b)_{n} n!}  \tag{9}\\
& \text {. }{ }_{1} F_{1}\left[\begin{array}{l}
\lambda+n ; \frac{x z}{z-1} \\
a+n ;
\end{array}{ }_{1} F_{1}\left[\begin{array}{l}
\lambda+n ; \frac{y z}{z-1} \\
b+n ;
\end{array}\right]\left[\frac{y y z}{(1-z)^{2}}\right]^{n},|z|<1 .\right.
\end{align*}
$$

Next we notice that the confluent hypergeometric function on the righthand side of (2) equals

$$
\begin{aligned}
& \sum_{n=0}^{\infty} \frac{(d)_{n}}{n!(1+a)_{n}}\left[\frac{x y z}{(1+z)^{2}}\right]^{n} \sum_{r, s=0}^{\infty} \frac{(d+n)_{r+s}(a-b)_{s}}{r!s!(1+a+n)_{s}}(-)^{r}\left(\frac{x z}{1-z}\right)^{r+s} \\
& =\sum_{n=0}^{\infty} \frac{(d)_{n}}{n!(1+a)_{n}}\left[\frac{x y z}{(1-z)^{2}}\right]^{n} \sum_{r=0}^{\infty} \frac{(d+n)_{r}}{r!}\left(\frac{x z}{z-1}\right)^{r} \sum_{s=0}^{r} \frac{(-r)_{s}(a-b)_{s}}{s!(1+a+n)_{s}}, \\
& =\sum_{n=0}^{\infty} \frac{(d)_{n}}{n!(1+a)_{n}}\left[\frac{x y z}{(1-z)^{2}}\right]^{n} \sum_{r=0}^{\infty} \frac{(d-n+n)_{r}(1+b+n)_{r}}{r!(1+a+n)_{r}}\left(\frac{x z}{z-1}\right)^{r},
\end{aligned}
$$

by Vandermonde's theorem [1, p. 3], and the bilateral generating relation (2) or (5) can thus be put in the equivalent form

$$
\begin{align*}
& \sum_{n=0}^{\infty} \frac{(b)_{n}}{n!}{ }_{2} F_{1}\left[\begin{array}{c}
-n, d ; \\
a ;
\end{array}\right]{ }_{1} F_{\perp}\left[\begin{array}{c}
-n ; y \\
b ;
\end{array}\right] z^{n} \tag{10}
\end{align*}
$$

Since

$$
\exp \left(\frac{y z}{z-1}\right)={ }_{1} F_{1}\left[\begin{array}{l}
b+n ; \cdots z \\
b+n ;
\end{array}\right],
$$

it is easy to see that (10) follows rather immediately from Meixner's formula ( [8], see also [5], p. 84)
*See also the formula (27 on page 288 of Erdelyi et al, [5]. For an elementary derivation of (9) with a slightly different right-hand side see the formula (2-3) of Srivastava [12].

$$
\begin{align*}
& \sum_{n=0}^{\infty} \frac{(\lambda)_{n}}{n!}{ }_{2} F_{1}\left[\begin{array}{c}
-n, d ; \\
a ;
\end{array}\right]{ }_{2} F_{1}\left[\begin{array}{c}
-n, e ; \\
b ;
\end{array}\right] z^{n}  \tag{11}\\
& =(1-z)^{-\lambda} \sum_{n=0}^{\infty} \frac{(\lambda)_{n}(d)_{n}(e)_{n}}{n!(a)_{n}(b)_{n}}{ }_{2} F_{1}\left[\begin{array}{r}
\left.\lambda+n, d+n ; \frac{x z}{z-1}\right] \\
a+n ;
\end{array}\right] \\
& \cdot{ }_{2} F_{1}\left[\begin{array}{r}
\left.\lambda+n, e+n ; \frac{y z}{z-1}\right]\left[\frac{x y z}{(1-z)^{2}}\right]^{n},|z|<1,
\end{array},\right.
\end{align*}
$$

when $y$ is replaced by $y / e, \lambda=b$, and $e \rightarrow \infty$.
The foregoing analysis when applied to the remaining formulas (8) and (1) will readily reduce them to their equivalent forms

$$
\begin{gather*}
\sum_{n=0}^{\infty} \frac{(b)_{n}}{n!}{ }_{2} F_{1}\left[\begin{array}{c}
-n, d ; \\
a ;
\end{array}\right]{ }_{1} F_{1}\left[\begin{array}{c}
-n ; y \\
1-b-n ;
\end{array}\right] z^{n}  \tag{12}\\
=(1-z)^{-b} \exp (y z) \sum_{n=0}^{\infty}-\frac{(d)_{n}}{(a)_{n}} \frac{(-x y z)^{n}}{n!}{ }_{2} F_{1}\left[\begin{array}{c}
b, d+n ; \frac{x z}{a+n ;} \overline{z-1}
\end{array}\right]
\end{gather*}
$$

and

$$
\begin{align*}
& \sum_{n=0}^{\infty} \frac{(b)_{n}}{n!}{ }_{1} F_{1}\left[\begin{array}{c}
-n ; \\
a ;
\end{array}\right]{ }_{2} F_{1}\left[\begin{array}{cc}
-n, d & ; \\
1-b-n ;
\end{array}\right] z^{n}  \tag{1.3}\\
& =(1-z)^{-b}{\underset{n}{n}=0}_{\infty}^{(a)} \underset{(a)}{(a)} \frac{(-x y z)^{n}}{n!}{ }_{1} F_{1}\left[\begin{array}{c}
b ; \\
a+n ; \\
z-1
\end{array}\right]{ }_{1} F_{0}\left[\begin{array}{c}
d+n: y z \\
-;
\end{array}\right]
\end{align*}
$$

respectively.
Both (12) and (13) are evidently immediate consequences of our bilinear generating relation (3.3) in [12], viz.

$$
\begin{align*}
& \sum_{n=0}^{\infty} \frac{(\lambda)_{n}}{n!}{ }_{p+1} F_{q}\left[\begin{array}{r}
-n, a_{1}, \ldots, a_{p} ; \\
b_{1}, \ldots, b_{q} ;
\end{array}\right] r_{r+1} F_{s+1}\left[\begin{array}{c}
-n, \alpha_{1}, \ldots, a_{r} ; \\
1-\lambda-n, \beta_{1}, \ldots, \beta_{s} ; y
\end{array}\right] z^{n} \tag{14}
\end{align*}
$$

$$
\begin{aligned}
& \text {. } p{ }_{+1} F_{q}\left[\begin{array}{c}
\lambda, a_{1}+n, \ldots, a_{p}+n ; \\
b_{1}+n, \ldots, b_{q}+n ;
\end{array} \frac{x z}{z-1}\right]{ }_{r} F_{s}\left[\begin{array}{c}
\alpha_{1}+n, \ldots, \alpha_{r}+n ; \\
\beta_{1}+n, \ldots, \beta_{s}+n ;
\end{array}\right], \\
& \text { [ } 411 \text { ] }
\end{aligned}
$$

valid for such values of $|x|,|y|$, and $|z|$ that the two sides have a meaning, provided that $p \leqslant q$ and $r \leqslant s+1$.

It may be of interest to remark that in the course of an attempt to give extensions of the well-known Hille-Hardy formulil (see, e.g., [6], p. 189)

$$
\begin{align*}
& \sum_{n=0}^{\infty} \frac{n!}{1^{\prime}(a+n+1)} \cdot L_{n}^{(a)}(x) L_{n}^{(a)}(y) z^{n}=(1-z)^{-1} \exp \left[-\frac{(x+y) z}{1-z}\right]  \tag{15}\\
& \quad \cdot(x y z)^{-\frac{1}{k} a} I_{a}\left[\begin{array}{c}
2 \sqrt{2}(x y z) \\
1-z
\end{array}\right], \quad|z|<1,
\end{align*}
$$

which follows readily from (9) when $a=b=\lambda$, we hive recently applied the Laplace and the inverse Laplace transform techniques and the method of multidimensional mithematical induction in order to obtain several new and distinct bilinear and bilateral generating functions for certain classes of generalized hypergeometric polynomials. These results include (14) and the bilinear generating function

$$
\begin{aligned}
& =(1-z)^{-\lambda} \sum_{n=0}^{\infty} \frac{(\lambda)_{n} \prod_{j=1}^{p}\left(a_{j}\right)_{n} \prod_{j=1}^{r}\left(n_{j}\right)_{n}}{n!\prod_{j=1}^{q}\left(b_{j}\right)_{n} \prod_{j=1}^{s}\left(\beta_{j}\right)_{n}}\left[\frac{x y z}{(1-z)^{n}}\right]^{n} \\
& \text { - }{ }_{p+1} F_{q}\left[\begin{array}{r}
\lambda+n, a_{1}+n, \ldots, a_{p}+n ; \frac{x z}{} \\
b_{1}+n, \ldots, b_{q}+n ; \overline{z-1}
\end{array}\right] r+1 F_{s}\left[\begin{array}{r}
\lambda+n, \alpha_{1}+n, \ldots,{ }_{q}+n ; y z \\
\beta_{1}+n, \ldots, \beta_{s}+n ; \frac{y z}{z-1}
\end{array}\right]
\end{aligned}
$$

which unify scores of hitherto scattered results in the theory of generalized hypergeometric polynomials (see [9], [10] and [15]). For instance, the special case $p=q=1$ of (16) corresponds to the main formula (4), p. 1320 of Brafman [2] when $\lambda-b_{1}$, and in the limit as $x \rightarrow 0$ our formulas (14) and (16), as also (2.8) in [12], reduce to the results proved carlier by Chaundy [3, p. 62]. For $p=q=r=s=1$, our bilinear relations (16), and (3•4) and (3.5) in [12], arc due to Meixner [8] who obtained them two decades ago by transforming the Pochhammer contour integral associated with Gauss's hypergeometic function.

We should like to conclude by mentioning that several extensions of the generating relations in [11] and [12] to hold for certain classes of general double hypergeometric functions appear in our subsequent works [13' and [14]. Indeed in terms of the Appell function [5, p. 224]

$$
F_{1}\left[\alpha, \beta, \beta^{\prime} ; \gamma ; x, \gamma\right]=\sum_{m, n=0}^{\infty} \frac{(\alpha)_{m+n}(\beta)_{m}\left(\beta^{\prime}\right)_{n}}{(\gamma)_{m+n}} \frac{x^{m}}{m!\frac{y^{n}}{n!}}
$$

we find it worthwhile to record the following very specialized, but nevertheless interesting, forms* of some of our results :

$$
\begin{align*}
& \sum_{n=0}^{\infty} \frac{(\lambda)_{n}}{n!}{ }_{2} F_{1}\left[\begin{array}{c}
n, a ; \\
b ;
\end{array}\right]{ }_{2} F_{1}\left[\begin{array}{cc}
-n, & \mu ; y \\
1-\lambda-n ;
\end{array}\right] z^{n}  \tag{17}\\
&=(1-z)^{-\lambda}(1-y z)^{-\mu} F_{1}\left[a, \lambda, \mu ; b ; \frac{x z}{z-1}, \frac{x y z}{y z-1}\right] ; \\
& \sum_{n=0}^{\infty} \frac{(\lambda)_{n}}{n!}{ }_{2} F_{1}\left[\begin{array}{c}
\rho-n, d ; \\
a ;
\end{array}\right]{ }_{2} F_{1}\left[\begin{array}{c}
\sigma-n, e ; \\
b ;
\end{array}\right] z^{n}  \tag{18}\\
&=(1-z)^{-\lambda} \sum_{n=0}^{\infty} \frac{(\lambda)_{n}(d)_{n}(e)_{n}}{n!(a)_{n}(b)_{n}} F_{1}\left[d+n, \rho, \lambda+n ; a+n ; x, \frac{x z}{z-1}\right] \\
& \quad F_{1}\left[e+n, \sigma, \lambda+n ; b+n ; y, \frac{y z}{z-1}\right]\left[\frac{x y z}{(1-z)^{2}}\right]^{n} .
\end{align*}
$$

The last formula reduces, when $\rho=\sigma=0$, to Meixner's bilinear generating function (1i) which, in turn, yields several hitherto scattered results including, e.g., $(4 \cdot 3)$ and (4.6) of Weisner [16, p-1037].
*Note that in view of (5) and (6), the results (17) and (18) do incorporate the four formulas (1) through (4) as their specialized or limiting cases.

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# Reactions of Ortho-esters of Germanium-reactions of ethyl Orthogermanate with acetic and Benzoic acids 

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Reactions of ethyl orthogermanate with acetic and benzoic acids in different stoichiometric ratios have been studied.

In a number of recent publications from these laboratories, interesting and useful results have been described in the study of the reactions of the alkoxy derivatives of aluminium ${ }^{1,2}$, silicon ${ }^{3-5}$, titanium ${ }^{6-8}$, zirconium ${ }^{9}, 10$ with carboxylic acids and anhydrides, (e.g., acetic acid, acetic anhydride, benzoic, lactic, mandelic and salicylic acids) in non-aqueous medium. In view of the above, it was considered of interest to study the reactions of ethyl orthogermanate with acetic and benzoic acids:

Ethyl orthogermanate on being allowed to react with one mole of acetic acid in refluxing benzene gave the triethoxy monoacetate derivative in almost quantitative yield:

$$
\mathrm{Ge}\left(\mathrm{OC}_{2} \mathrm{H}_{5}\right)_{4}+\mathrm{CH}_{3} \mathrm{COOH} \rightarrow\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{3} \mathrm{Ge}\left(\mathrm{OOC}^{-\mathrm{CH}_{3}}\right)+\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{OH}
$$

when the above reaction was carried out in $1: 2$ molar ratio, only 1.7 moles of alcohol could be fractionated out azeotropically and a basic acetate, $\mathrm{GeO}_{1} \cdot \delta$ $\left(\mathrm{OOCCH}_{3}\right)_{0 \cdot 25}\left(\mathrm{OC}_{2} \mathrm{H}_{5}\right)_{0.78}$ was isolated from the reaction mixture.

In order to avoid the type of decomposition observed above, the reactions of ethyl orthogermanate with acetic acid were carried out in cold in 1:2 and 1:4 molar ratios when the resulting products were found to have acetate : germanium ratio of $0.88: 1$ and $1.08: 1$ respectively. The former compound on being heated in a bath of $100-110^{\circ}$ for about half an hour appeared to yield a derivative with acetate : germanium ratio of $0.66: 1$.

The reaction between ethyl orthogermanate and acetic anhydride ( $1: 4$ ) in refluxing benzene yielded a product which was found to have acetate: germanium ratio of $0.89: 1$.

It appears that the main straight-forward reaction of ethyl orthogermanate with acetic acid or anhydried can be represented by the following equations :

$$
\begin{aligned}
& \mathrm{Ge}\left(\mathrm{OC}_{2} \mathrm{H}_{5}\right)_{4}+\mathrm{CH}_{3} \mathrm{COOH} \rightarrow\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{3} \mathrm{Ge}\left(\mathrm{OOCCH}_{3}\right)+\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{OH} \\
& \left.\mathrm{Ge}\left(\mathrm{OC}_{2} \mathrm{H}_{5}\right)_{4}+\mathrm{CH}_{3} \mathrm{CO}\right)_{2} \mathrm{O} \rightarrow\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{3} \mathrm{Ge}\left(\mathrm{OOCCH}_{3}\right)+\mathrm{CH}_{3} \mathrm{COOC}_{2} \mathrm{H}_{5} .
\end{aligned}
$$

Further reactions are slow and accompanied with side decomposition.
It is interesting to recall that in the corresponding reactions of alkyl orthotitanates ${ }^{8}$ with acetic acid or anhydried, a diacetate derivative is formed by an
exothermic facile reaction and in this case also further reaction is much slower and accompanied by side reaction which mainly consist in the elimination of ethyl acetate :

$$
\begin{aligned}
& \mathrm{Ti}\left(\mathrm{OC}_{2} \mathrm{H}_{6}\right)_{4}+2\left(\mathrm{CH}_{3} \mathrm{CO}\right)_{2} \mathrm{O} \xrightarrow[\rightarrow]{\text { fast }}\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{2} \mathrm{Ti}\left(\mathrm{OOCCH}_{3}\right)_{2}+2 \mathrm{CH}_{3} \mathrm{COOC}_{2} \mathrm{H}_{5} \\
& \left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{2} \mathrm{Ti}\left(\mathrm{OOCCH}_{3}\right)_{2}+\left(\mathrm{CH}_{3} \mathrm{CO}\right)_{2} \mathrm{O} \xrightarrow{\text { sow }}\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right) \mathrm{Ti}\left(\mathrm{OOCCH}_{3}\right)_{3}+ \\
& \left(\mathrm{C}_{2} \mathrm{H}_{6} \mathrm{O}\right) \mathrm{Ti}\left(\mathrm{OOCCH}_{3}\right)_{3} \rightarrow \mathrm{O}=\mathrm{Ti}\left(\mathrm{OOCCH} \mathrm{O}_{3}\right)_{2}+\mathrm{CH}_{3} \mathrm{COOC}_{2} \mathrm{H}_{5}
\end{aligned}
$$

In view of the above interesting results the reaction of germanium tetracetate with ethanol ( $1: 4$ ) was also carried out at room temperature and the resulting product was shown to have acetate : germanium ratio of $1 \cdot 05: 1$.

$$
\mathrm{Ge}\left(\mathrm{OOCCH}_{3}\right)_{4}+3 \mathrm{C}_{2} \mathrm{H}_{5} \mathrm{OH} \rightarrow\left(\mathrm{C}_{2} \mathrm{H}_{6} \mathrm{O}\right)_{3} \mathrm{Ge}\left(\mathrm{OOCCH}_{3}\right)+3 \mathrm{CH}_{3} \mathrm{COOH}(4 \text { moles })
$$

The reactions between ethyl orthngermanate and benzoic acid have also been studied in different stoichiometric ratios. The reactions follow a pattern quite similar to the course of reactions, described by Mehrotra and coworkers for alkyl ortho-titanate :

$$
\begin{align*}
& \mathrm{Ge}\left(\mathrm{OC}_{2} \mathrm{H}_{5}\right)_{4}+\mathrm{C}_{6} \mathrm{H}_{5} \mathrm{COOH} \rightarrow\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{3} \mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{5}\right)+\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{OH}  \tag{l}\\
& \mathrm{I} \\
& \mathrm{Ge}\left(\mathrm{OC}_{2} \mathrm{H}_{5}\right)_{4}+2 \mathrm{C}_{6} \mathrm{H}_{5} \mathrm{COOH} \rightarrow\left(\mathrm{C}_{2} \mathrm{H}_{8} \mathrm{O}\right)_{2} \mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{8}\right)_{2}+2 \mathrm{C}_{2} \mathrm{H}_{8} \mathrm{OH}  \tag{2}\\
& \mathrm{II} \\
& \mathrm{Ge}\left(\mathrm{OC}_{3} \mathrm{H}_{5}\right)_{4}+3 \mathrm{C}_{6} \mathrm{H}_{5} \mathrm{COOH} \rightarrow\left(\mathrm{C}_{2} \mathrm{H}_{8} \mathrm{O}\right) \mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{5}\right)_{3}+3 \mathrm{C}_{2} \mathrm{H}_{5} \mathrm{OH}  \tag{3}\\
& \text { III }
\end{align*}
$$

The preparation of germanium tetrabenzoate was also attempted by refuxing ethyl orthogermanate during the course of which only 2.87 moles of alcohol could be fractionated out azcotropically. On allowing the reaction mixture to reflux for a few more hours, a fine white solid separated. This was heated under reduced pressure to give a light amorphous powder which corresponded in analysis to germanyl benzoate. The reactions can, thercfore, be represented by the following equations:

$$
\begin{align*}
& \mathrm{Ge}\left(\mathrm{OC}_{42} \mathrm{H}_{5}\right)_{4}+3 \mathrm{C}_{6} \mathrm{H}_{5} \mathrm{COOH} \rightarrow\left(\mathrm{C}_{2} \mathrm{H}_{6} \mathrm{O}\right) \mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{5}\right)_{3}+3 \mathrm{C}_{2} \mathrm{H}_{8} \mathrm{OH} \\
& \left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right) \mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{5}\right)_{3} \rightarrow\left(\mathrm{C}_{6} \mathrm{H}_{5} \mathrm{COO}\right)_{2} \mathrm{Gc}=\mathrm{O}+\mathrm{C}_{6} \mathrm{H}_{5} \mathrm{COOC}_{2} \mathrm{H}_{5} \tag{4}
\end{align*}
$$

The formation of germanyl benzoate in the reaction (4) was confirmed by heating the compound (III) under reduced pressure :

$$
\left(\mathrm{C}_{2} \mathrm{H}_{8} \mathrm{O}\right) \mathrm{Ge}\left(\mathrm{OOCC}_{5} \mathrm{H}_{8}\right)_{3} \xrightarrow[\text { red. prcss. }]{\text { heat }} \mathrm{O}=\mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{5}\right)_{2}+\mathrm{C}_{8} \mathrm{H}_{5} \mathrm{COOC}_{2} \mathrm{H}_{5}
$$

Diethoxy germanium dibenzoate, described above (II), was found to interchange its ethoxy groups with higher alcohols, e.g. n-butanol and tertiary butanol:

$$
\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{2} \mathrm{Ge}\left(\mathrm{OOCC}_{6}\left(\mathrm{H}_{5}\right)_{2}+2 \mathrm{R}^{\prime} \mathrm{OH} \rightarrow\left(\mathrm{R}^{\prime} \mathrm{O}\right)_{2} \mathrm{Gc}\left(\mathrm{OOCG}_{6} \mathrm{H}_{8}\right)_{2}+2 \mathrm{C}_{2} \mathrm{H}_{5} \mathrm{OH}\right.
$$

The alcohol produced in all the above reactions was removed azeotropically with benzene, and the course of reaction was followed both by estimating the liberated alcohol in the azeotrope and the analysis of germanium in the products.

## Experimental

All glass apparatus with interchangeable joints was used and moisture was excluded from all the reactions.

Analytical procedures and methods for drying the reagents have been described elsewhere ${ }^{11}$. Ethyl orthogermanate was prepared by the ammonia method ${ }^{12}$. Germanium tetraacetate ${ }^{13}$ was prepared by reacting germanium tetrachloride with acetic anhydride in the presence of thallous acetate. Acetic acid and anhydride were purified by fractionation. Benzoic acid (B. D. H., analar) was dried at $40^{\circ} / 15 \mathrm{~mm}$ immediately before use.

1. Reaction between ethyl orthogermanate and acetic acid; molar ratio 1:1.

Acetic acid ( 0.38 g ) was admitted to a mixture of ethyl orthogermanate ( 1.58 g ) and benzene $(40 \mathrm{~g})$. On shaking the reaction mixture, a clear solution was obtained which was refluxed under the fractionating column for about two hours at $100^{\circ}$. The ethanol liberated in the reaction was fractionated off azeotropically. The excess of the solvent was removed under reduced pressure and the compound was dried at $32^{\circ} / 0.5 \mathrm{~mm}$. for about two hours. A colourless liquid ( 1.60 g ), slightly more viscous than ethyl orthogermanate, fuming in moist air, was obtained.

Found : ethanol in the azeotrope, 0.26 g ( 1 mole requires 0.28 g ).
$\%$, found : $\mathrm{Ge}, 27 \cdot 20 ; \mathrm{OC}_{2} \mathrm{H}_{5}, 50 \cdot 63$; $\mathrm{OOCCH}_{3}, 22 \cdot 13$; Calc. for
$\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{3} \cdot \mathrm{Ge}(\mathrm{OOCGH})_{3}: \mathrm{Ge}, 27 \cdot 21 ; \mathrm{OC}_{2} \mathrm{H}_{5}, 50 \cdot 66 ; \mathrm{OOCGH}_{3}, 22 \cdot 13$.
2. Reaction between ethpl orthogermanate and acetic acid; molar ratio $1: 2$.

Acetic acid ( $1 \cdot 14 \mathrm{~g}$ ) was treated with ethyl orthogermanate $(2.82 \mathrm{~g})$ in refluxing benzene $(40 \mathrm{~g})$ for about three hours, and the benzene-ethanol azeotrope was collected very slowly. More benzene ( 20 g ) was added and the mixture was refluxed for another two hours and the collection of the azeotrope was continued. The remaining solution was evaporated to dryness under reduced pressure. A colourless semi-solid miscible with benzene together with some insoluble impurity was obtained. The soluble portion was extracted with benzene. On drying it under reduced pressure a spongy white solid $(1 \cdot 09 \mathrm{~g})$, soluble in benzene, was obtained.

Fourd : ethanol in the azeotrope, 0.89 g ( 2 moles require 1.027 g ).
$\%$, found : Ge, 49.91 ; $\mathrm{OC}_{2} \mathrm{H}_{5}, 23.20$; $\mathrm{OOCCH}_{3}, 9.93$; Calc. for
$\mathrm{GeO}_{1 \cdot 5}\left(\mathrm{OC}_{2} \mathrm{H}_{5}\right)_{0 \cdot 75}\left(\mathrm{OOCCH}_{3}\right)_{0 \cdot 25}: \mathrm{Ge}, 50 \cdot 01 ; \mathrm{OC}_{2} \mathrm{H}_{5}, 23 \cdot 28 ; \mathrm{OOCCH}_{3}, 10 \cdot 17$.
3. Reaction between ethyl orthogermanate and acetic acid in cold; molar ratio $1: 2$.

Acetic acid $(0.99 \mathrm{~g})$ was introduced to a mixture of ethyl orthogermanate $(2.06 \mathrm{~g})$ and benzene $(2 \mathrm{~g})$. On shaking the reactants for some time a slight turbidity appeared which, however, redissolved on keeping the reaction mixture overnight. The solvent was removed under reduced pressure and the product was dried at $32^{\circ} / 0.5 \mathrm{~mm}$. for about two hours. A colourless liquid ( 1.91 g ) was obtained.
$\%$, found : $\mathrm{Ge}, 28.40 ; \mathrm{OC}_{2} \mathrm{H}_{5}, 50 \cdot 40 ; \mathrm{OOCCH}_{3} ; 20 \cdot 50$.
(Ge $: \mathrm{OC}_{2} \mathrm{H}_{5}$ ratio $=1: 2 \cdot 86 ; \mathrm{Ge}: \mathrm{OOCGH}_{3}$ ratio $=1: 0.88$ ).
The above compound $(1.54 \mathrm{~g})$ was heated in a bath of $100-110^{\circ}$ for about 30 minutes and then evacuated under reduced pressure for about two hours. A colourless liquid $(1 \cdot 27 \mathrm{~g})$ was obtained.
$\%$, found : $\mathrm{Ge}, 30.58 ; \mathrm{OC}_{2} \mathrm{H}_{5}, 49.00 ; \mathrm{OOCCH}_{3}, 16.50$.
(Ge : $\mathrm{OC}_{2} \mathrm{H}_{5}$ ratio $=1: 2 \cdot 58 ; \mathrm{Ge}: \mathrm{OOCCH}_{3}$ ratio $=1: 0.66$ ).
4. Reaction between ethyl orthogermanate and acetic acid in cold ; molar ratio $1: 4$.

A slight turbidity appeared on shaking a mixture of ethly orthogermanate $(1.65 \mathrm{~g})$, acetic acid ( $1 \cdot 57 \mathrm{~g}$ ) and benzene ( 2 g ). The mixture was kept overnight and then evacuated under reduced pressure at $34^{\circ} / 0.5 \mathrm{~mm}$. for about three hours. A colourless, fuming liquid $(1.52 \mathrm{~g})$ was obtained.
$\%$, found : Ge, $29.00 ; \mathrm{OC}_{2} \mathrm{H}_{5}, 48 \cdot 00 ; \mathrm{OOCCH}_{3} ; 25 \cdot 50$.
(Ge: $\mathrm{OC}_{3} \mathrm{H}_{5}$ ratio $=1: 2 \cdot 67 ; \mathrm{Ge}: \mathrm{OOCGH}_{3}$ ratio $=1: 1 \cdot 08$ ).
5. Reaction between ethyl orthogermanate and acetic anhydride; molar ratio $1: 4$.

Acetic anhydride $(2 \cdot 80 \mathrm{~g})$ was added dropwise to a solution of ethyl orthogermanate $(1.73 \mathrm{~g})$ in benzene $(40 \mathrm{~g})$. On shaking a slight turbidity was noticed which disappeared on keeping the reactants overnight. The reaction mixture was refluxed at about $100^{\circ}$ for $3-4$ hours. The bath temperature was then raised and the distillate coming at $80^{\circ}$ was collected. About 25 g . of the distillate was withdrawn and the rest of the solvent was removed under reduced pressure and the product was dried at $55-60^{\circ} / 0 \cdot 5-1 \cdot 0 \mathrm{~mm}$ for about two hours. A slight turbid liquid $(1 \cdot 40 \mathrm{~g})$ was obtained.
$\%$, found : Ge, $28 \cdot 60 ; \mathrm{OC}_{2} \mathrm{H}_{5}, 49 \cdot 60 ; \mathrm{OOCCH}_{3}, 20.80$.
(Ge : $\mathrm{OC}_{2} \mathrm{H}_{5}$ ratio $=1: 2.8 ; \mathrm{Ge}: \mathrm{OOCCH}_{3}$ ratio $=1: 0.89$ ).
6. Reaction between germanium tetra-acetate and ethanol in cold; molar ratio 1:4.

Noticeable heat was produced when a mixture of germanium tetra-acetate $(1.06 \mathrm{~g})$, ethanol $(0.63 \mathrm{~g})$ and benzene ( 20 g ) was shaken together. The excess of the solvent was removed under reduced pressure and the compound was dried at $32^{\circ} / 1 \mathrm{~mm}$. for about two hours. A colourless moblile liquid ( 0.76 g ) was obtained.
$\%$, found: $\mathrm{Ge}, 26.00 ; \mathrm{OC}_{2} \mathrm{H}_{5}, 46.00 ; \mathrm{OOCCH}_{3}, 22 \cdot 20$.
( $\mathrm{Ge}: \mathrm{OC}_{2} \mathrm{H}_{5}$ ratio $=1: 2.85 ; \mathrm{Ge}: \mathrm{OOCCH}_{3}$ ratio $=1: 1 \cdot 05$ ).
Thus in experiments (3-6) the final product was mainly, $\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{3} \mathrm{Ge}\left(\mathrm{OOGCH}_{3}\right)$, which requires $\mathrm{Ge}, 27 \cdot 21 ; \mathrm{OC}_{2} \mathrm{H}_{5}, 50 \cdot 66$; $\left(\mathrm{OOCGH}_{3}, 22 \cdot 13 \%\right.$.
7. Reaction between ethyl orthogermanate and benzoic acid; malar ratio 1:1.

Benzoic acid $(0.97 \mathrm{~g})$, ethyl orthogermanate $(2.02 \mathrm{~g})$ and benzene ( 40 g ) were shaken together. The clear solution, thus obtained, was refluxed under the fractionating column at $100-150^{\circ}$ for about 1.5 hours. The ethanol liberated in the reaction was completely fractionated off azeotropically. The reaction maxture was rendered free of solvent under reduced pressure and the residual product was dried at $50^{\circ} / 0.5 \mathrm{~mm}$. for 1.5 hours. A whitish slightly viscous liquid $(2.58 \mathrm{~g})$ was obtained.

Found : ethanol in the azeotrope; 0.36 g ( 1 mole requires 0.368 g ).
$\%$, found : $\mathrm{Ge}, 21 \cdot 90$; Calc. for $\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{3} \mathrm{Ge}\left(\mathrm{OOCG}_{8} \mathrm{H}_{5}\right): \mathrm{Ge}, 22 \cdot 07$.
On heating 1.5 g of the compound under reduced pressure, a colourless liquid ( 0.8 g ) b.p. $75-80^{\circ} / 3-4 \mathrm{~mm}$. was obtained.
$\%$, found: $\mathrm{Ge}, 15 \cdot 10$.
The residue in the flask contained $\mathrm{Ge}, 43 \cdot 64 \%$.
8. Reaction between ethyl orthogermanate and benzoic acid; molar.ratio 1:2.

Benzoic acid $(1.87 \mathrm{~g})$ was admitted to a solution of ethyl orthogermanate $(1 \cdot 94 \mathrm{~g})$ in benzene $(45 \mathrm{~g})$. On shaking a colourless solution was obtained which was allowed to reflux for about two hours at $110^{\circ}$. The distillate was withdrawn dropwise between $68-80^{\circ}$. The pure benzene (b. p. $80^{\circ}$ ) was distilled under high reflux ratio. The product was freed of solvent under reduced pressure at $50^{\circ} / 1 \mathrm{~mm}$. A whitish viscous liquid ( $2 \cdot 99 \mathrm{~g}$ ), miscible with benzene, was obtained.

Found ; ethanol in the azeotrope, $0.70 \mathrm{~g}(2$ moles require 0.70 g$)$.
$\%$, found : $\mathrm{Ge}, 17.90$; Calc. for $\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{2} \mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{5}\right)_{2}$ : $\mathrm{Ge}, 17.92$.
9. Reaction between $\left(\mathrm{C}_{3} \mathrm{H}_{5} \mathrm{O}\right)_{2} \mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{5}\right)_{2}$ and excess normal butanol :

The compound, (from experiment, 8$),(1.80 \mathrm{~g})$, normal butanol $(3.0 \mathrm{~g})$ and benzene ( 30 g ), were refluxed for about two hours at $110^{\circ}$. The ethanol liberated in the reaction was fractionated out very slowly. The remaining solvent and butanol were removed under reduced pressure and the compound was dried at $40^{\circ} / 1 \mathrm{~mm}$. for about two hours. A white semi-solid $(2 \cdot 00 \mathrm{~g})$, soluble in benzene, was obtained.
$\%$, found : Ge, 15.69 ; Calc. for $\left(\mathrm{C}_{4} \mathrm{H}_{8} \mathrm{O}^{n}\right)_{2} \mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{5}\right)_{2}$ : Ge, 15.74 .
10. Reaction between $\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right)_{2} \mathrm{Ge}\left(\mathrm{OOCC}_{2} \mathrm{H}_{5}\right)_{2}$ and excess tert. butanol :

The compound (from experiment, 2), ( 1.82 g ), tertiary butanol ( $4 \cdot 45 \mathrm{~g}$ ) and benzene ( 40 g ), were allowed to reflux at $115-20^{\circ}$ as usual. The ethanol-benzene azeotrope was slowly collected. The volatile fractions were evaporated under reduced pressure and the compound was dried at $45^{\circ} / 2 \mathrm{~mm}$. for about two hours. A white semi-solid ( $1 \cdot 98 \mathrm{~g}$ ) miscible with benzene, was obtained.

$$
\text { Found : eihanol in the azeotrope, } 0.38 \mathrm{~g} \text { ( } 2 \text { moles require } 0.4 \mathrm{lg} \text { ). }
$$

$$
\% \text { found : Ge, } 15 \cdot 80 \text {; Calc. for }\left(\mathrm{C}_{4} \mathrm{H}_{9} t \mathrm{O}\right)_{2} \mathrm{Ge}\left(\mathrm{OOC}_{6} \mathrm{H}_{5}\right)_{2}: \mathrm{Ge}, 15 \cdot 74
$$

11. Reaction between ethyl orthogermanate and benzoic acid; molar ratio $1: 3$.

Benzoic acid ( 3.46 g ) was added with the help of benzene ( 50 g ) to ethly orthogermanate $(2.39 \mathrm{~g})$ and the reaction mixture was then refluxed for about three hours at a bath temperature of $120^{\circ}$. The distillate was collected dropwise at $68^{\circ}$. The temperature of the distillate then rose to $76^{\circ}$ after nearly 8 cc . of it were withdrawn. At this stage the reaction misture was allowed to reflux more, for about two hours. The azeotrope was again collected dropwise till the temperature of distillate became steady ai $80^{\circ}$ and showed no tendency to decrease. The subsequent fractions of the disillate were collected separately under a high reflux ratio, 1:20. The product was freed of solvent under reduced pressure at $52^{\circ} / 1 \mathrm{~mm}$. A white semi-solid ( $4 \cdot 3 \mathrm{lg}$ ), soluble in benzene was obtained.

Found : ethanol in the azeotrope (fractions collected other than those distilling at $80^{\circ}$ ), 1.25 g ( 3 moles require $1 \cdot 30 \mathrm{~g}$ ).
$\%$, found : $\mathrm{Ge}, 15 \cdot 17$; Calc. for $\left(\mathrm{C}_{2} \mathrm{H}_{5} \mathrm{O}\right) \mathrm{Ge}\left(\mathrm{OOCG}_{6} \mathrm{H}_{5}\right)_{3}: \mathrm{Ge}, 15 \cdot 09$.
Above compound ( 2.0 g ) was heated under reduced pressure at $100-150^{\circ} / 0.05$ mm . for 2.5 hours. A white powder $(1.28 \mathrm{~g})$ insoluble in benzene, was obtained.
$\%$, found: $\mathrm{Ge}, 21 \cdot 01$; Calc. for $\mathrm{O}=\mathrm{Ge}\left(\mathrm{OOCC}_{6} \mathrm{H}_{5}\right)_{2} \mathrm{Ge}, 21 \cdot 93$.
12. Reaction between ethyl orthogermanate and excess ( $<4$ moles) benzoic acid :

Ethyl orthogermanate ( $2 \cdot 06 \mathrm{~g}$ ), benzoic acid $(4 \cdot 48 \mathrm{~g})$ and benzene ( 60 g ) were refluxed under the column for many hours ( 6 hours). The ethanol formed was
removed azeotropically. A white finely divided solid began to separate when nearly all the ethanol was fractionated off. At this stage, the reaction mixture was further refluxed for about eight hours when more solid separated. It was then allowed to stand overnight. The supernatant liquor was decanted off along with some solid mass. The residue was thoroughly washed with dry benzene and dricd at $70^{\circ} / 0 \cdot 1 \mathrm{~mm}$. A white powder $(2 \cdot 60 \mathrm{~g})$ was obtained.

Found : ethanol in the azeotrope, $1 \cdot 072 \mathrm{~g}$ ( 3 moles require $1 \cdot 12 \mathrm{~g}$ ).
The above solid ( $2 \cdot 0 \mathrm{lg}$ ), was heated in a bath of $105^{\circ} / 0.2 \mathrm{~mm}$. for about two hours when droplets of a colourless liquid (having fruity smell) appeared in the receiver. A white powder ( 1.7 Ig ), insoluble in benzene, was left in the distill. ing flask.
$\%$, found : $\mathrm{Ge}, 21 \cdot 09$; Calc. for $\left[\mathrm{O}=\mathrm{Ge}\left(\mathrm{OOCC}_{0} \mathrm{H}_{5}\right)_{2}\right]_{n}: \mathrm{Ge}, 21 \cdot 95$.

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# On Resultant Kernel Functions 

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#### Abstract

In this paper, some rules connecting different classes of self-reciprocal functions have been stated for the generalised Hankel transform, which may be introduced by using a symmetrical Fourier kernel due to Roop Narain. Further, by making use of these rules and Hardy's Formula to give resultant of two kernels $$
K(x)=\int_{0}^{\infty} P_{1}(x) P_{2}(x y) d y
$$ where $P_{1}(x)$ and $P_{2}(x)$ are Fourier kernels, certain theorems have been established to investigate the self-reciprocity of functions in the generalised Hankel transform. Various particular cases have also been discussed.


## 1. The Hankel transform

$$
g(x)=\int_{0}^{\infty}(x y)^{\frac{1}{2}} J_{\nu}(x y) f(y) d y
$$

may be generalised by using a symmetrical Fourier kernel, given by Narain[9] in the form :

$$
g(x)=2 \beta \gamma \int_{0}^{\infty}(x y)^{\gamma-\frac{1}{2}} G_{2 p, 2 q}^{q, p}\left[\begin{array}{l|l}
\beta^{2}(x y)^{2 \gamma} & \left.\left.\begin{array}{l}
a_{1}, \ldots, a_{p},-a_{1}, \ldots,-a_{p} \\
b_{1}, \ldots, b_{q},-b_{1}, \ldots,-b_{q}
\end{array}\right] f(y) d y\right) .
\end{array}\right]
$$

where $\beta$ and $\gamma$ are real constants.
With $\beta=\frac{1}{2}, \gamma=1, p=0, q=1$ and $b_{1}=\nu / 2,(1 \cdot 2)$ reduces to ( $1 \cdot 1$ ). Also, under certain substitutions [91, the kernel of (1.2) yields, as particular cases, various kernels, given by Watson [13], Bhatnagar [2], Narain \{ [6] and [7] \} and Everitt [5].

In $(1 \cdot 2)$, if $f(x) \equiv g(x)$, then $f(x)$ will be said to be 'self-reciprocal' i.e. $K\left(a_{p} ; b_{q}\right)$ in the generalised Hankel transform ( $1 \cdot 2$ ).

The object of this paper is firstly to state some rules connecting different classes of self-reciprocal functions in $(1 \cdot 2)$ and to use these rules in establishing certain theorems which are useful in identifying the nature of resultant of two kernels as transforming a self-reciprocal function of known order into another of a different order. Some of these types of kernels have been obtained by making use of the theorems and various known and unknown resuls, which follow as particular cases, have also been discussed.

In the present discussion, for the sake of brevity, the notation of the type ${ }_{m} \lambda_{n}$ has been used to stand for a set of parameters $\lambda_{m}, \lambda_{m+1}, \ldots, \lambda_{n}$.
2. Rule $I$ : If $f(x)$ is $R\left(a_{p} ; b_{q}\right)$ and belongs to $A(\alpha, a)$ [12, p. $252 \S 9 \cdot 7$ ], then

$$
g(x)=(1 / x) \int_{0}^{\infty} P(y / x) f(y) d y
$$

is $R\left(c_{p}, ; d_{q}\right)$, provided that

where $X(s)$ is regular and satisfies the condition

$$
\chi(s)=\chi(1-s), \quad s=\sigma+i t
$$

in the strip

$$
a<\sigma<1-a
$$

and

$$
x(s)=0\left({ }_{e}^{\left[\left(t q+q^{\prime}-p-p^{\prime}\right) \pi / \gamma-\alpha+\eta\right]|t|}\right)
$$

for every positive $\eta$ and uniformly in any strip interior to $(2 \cdot 2)$ and $c$ is any value of $\sigma$ in (2.2)

In other words, $P(x)$ is a kernel transforming an $R\left(a_{p} ; b_{q}\right)$ into an $R\left(c_{p} ; d_{q^{\prime}}\right)$.
Rule II: If $f(x)$ is $R\left(a_{p} ; b_{q}\right)$ and belongs to $A(\alpha, a)$ [12, p. $252 \S 9 \cdot 7$ ], then the function

$$
g(x)=\int_{0}^{\infty} P(x y) f(y) d y
$$

is $R\left(c_{p} ; ; d_{q^{\prime}}\right)$, provided that
 where $\chi(s)$ satisfies the conditions $(2 \cdot 1)$ and ( $2 \cdot 3$ ). Since the result is symmetrical, the kernel $(2 \cdot 4)$ also effects the converse transformation.

In other words, $P(x)$ is a kernel transforming an $R\left(a_{p} ; b_{q}\right)$ into an $R\left(c_{p} ; d_{q}\right)$ and vice versa.

The proofs of the above rules can be developed on the same lines as in the proofs of the corresponding theorems of the Hankel transform [12, pp. 268-270].
3. Theorem 1 : The resultant

$$
K(x)=\int_{0}^{\infty} P_{1}(y) P_{2}(x y) d y,
$$

of the kernels of the form :
and
where

$$
0<c<1,0<c^{\prime}<1 \text { and }
$$

$$
\chi(s)=\chi(1-s) ; \quad \omega(s)=\omega(1-s),
$$

is a kernel transforming an $R\left(c_{p_{1}} ; d_{q_{1}}\right)$ into an $R\left(\alpha_{p_{2}} ; \beta_{q_{2}}\right)$ and vice-versa.
$P$,oof : Putting for $P_{2}(x y)$ in (3.1) from (3.3) and then changing the order of integration, the resultant kernel becomes as

$$
\begin{aligned}
& \omega(s) x^{-s} d s \int_{0}^{\infty} y^{-8} P_{\mathrm{L}}(\gamma) d y .
\end{aligned}
$$

The value of $y$-integral in (3.5) ean easily be deduced from (3.2) by using Mellin's inv ersion formula and therefore substituting it in (3.5), we get

$$
K(x)=\frac{1}{2 \pi i} \int_{c^{\prime}-i_{\infty}}^{c^{\prime}+i \infty} \beta^{-s / \gamma} \frac{{\underset{j}{j=1}}_{q_{1}} \Gamma\left(\frac{2 \gamma-1}{4 \gamma}+d_{j}+\frac{s}{2 \gamma}\right) \prod_{j=1}^{p_{1}} \Gamma\left(\frac{2 \gamma-1}{4 \gamma}-c_{j}+\frac{s}{2 \gamma}\right){\underset{j}{j=1}}_{p_{j}} \Gamma\left(\frac{2 \gamma-1}{4 \gamma}+\beta_{j}+\frac{s}{2 \gamma}\right)}{\left.p_{j}-\alpha_{j}+\frac{s}{2 \gamma}\right)} \theta(s) x^{-s} d s
$$

where
which clearly satisfies the functional equation

$$
\theta(s)=\theta(1-s)
$$

Hence, by virtue of Rule II, given in section 2, it is obvious from the form of (3.6) that $K(x)$, defined by (3•l), is a kernel which transforms an $R\left(\alpha_{p_{2}} ; \beta_{q_{2}}\right)$ into an $R\left(c_{p_{1}} ; d_{q_{2}}\right)$ and vice-versa.

It only remains to justify the change in the order of integration in (3.5). Now, putting $s=c^{\prime}+i t$ and using

$$
\left|\Gamma\left(A+\frac{1}{2} i t / \gamma\right)\right| \approx\left(e^{-\frac{\pi}{4 \gamma}|t|}|t|^{A-\frac{1}{2}}\right) \text {, as }|t| \rightarrow \infty
$$

in the integrand of (3.6), we see that its modulus is less than a constant multiple of

$$
\begin{align*}
& \left(\beta^{\gamma} x\right)^{-0 \prime} \int_{-\infty}^{\infty} e^{\infty}-\frac{\pi}{4 \gamma}\left(q+q_{2}-p-p_{2}\right)|t||t| \sum_{1}^{p} a_{j}+\sum_{1_{j}}^{q} b_{j}+\sum_{1}^{p_{2}} \alpha_{j}+{\underset{1}{\Sigma} \beta_{j}}_{q_{2}} \\
& \quad+\left(\frac{2 c^{\prime}-1}{4 \gamma}\right)\left(q+q_{2}-p-p_{2}\right)\left|\omega\left(c^{\prime}+i t\right)\right| d t \int_{0}^{\infty}\left|P_{1}(y)\right| y^{-c^{\prime}} d y .
\end{align*}
$$

The change of the order of integration will be permissible if both the $y$ and $t$ - integrals in (3.7) exist. These will exist if both the functions $P_{1}(x)$ and $P_{2}(x)$ belong to $A(\alpha, a)[12$, p. 252 § 9.7]

Now, the modulus of the integrand of (3.3) is less than a constant multiple of

$$
\begin{aligned}
& r-\sigma e^{-(\alpha-\eta-|\theta|) \mid t}|t|_{1}^{p} \alpha_{j}+\sum_{1}^{* q} b_{j}+\sum_{1}^{p_{2}} \alpha_{j}+\sum_{i}^{q_{2}} \beta_{j}+\left(\frac{2 \sigma-1}{4 \sigma}\right)\left(q+q_{2}-p-p_{2}\right), \\
& \text { since }
\end{aligned}
$$

$$
\omega\left(c^{\prime}+i t\right)=0\left(e^{\frac{\pi}{4 \gamma}\left(q+q_{2}-p-p_{2}-\alpha+\eta\right)|t|}\right)
$$

for every positive $\eta$ and uniformly in any strip interior to $a<\sigma<1-a$, where $a^{\prime}$ is any value of $\sigma$ in this strip. Thus the integral converges uniformly in any domain of $x$ interior to the angle $\alpha$ and in any smaller angle it is $0\left(|x|^{-c}\right)$.

Since $c^{\prime}$ can be taken as near to $(1-a)$ or $a$ as we please, $P_{2}(x)$ belongs to $A(x, a)$. Similarly, $P_{1}(x)$ can be shown to belong to $A(a, a)$.

This completes the proof of the theorem 1 .
3.1. Corollaries: (i) With $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=1, q_{1}=q_{2}=q=2$, $a_{1}=k^{\prime}-m^{\prime}-\frac{1}{3}-v / 2, \quad b_{1}=v / 2, \quad b_{2}=\nu / 2+2 m^{\prime}, \quad c_{1}=k-m-\frac{1}{2}-\mu / 2, \quad d_{1}=\mu / 2$, $d_{3}=\mu / 2+2 m, \alpha_{1}=l-n-\frac{1}{2}-\rho / 2, \beta_{1}=\rho / 2$ and $\beta_{2}=\rho / 2+2 n$, the above theorem reduces to a result, given by Saxena [10, p. 88] for $\chi_{\nu, k, m}$-transform.
(ii) On having $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=0, q_{1}=q_{2}=q=1, b_{1}=v / 2$, $d_{1}=\mu / 2$ and $\beta_{1}=\rho / 2$, the theorem yields a known result due to Agrawal [1, p. 305], in the Hankel transform.
3.2. Example: Let us choose the kernels [11]

$$
\begin{array}{r}
P_{1}(x)=x^{\gamma-\frac{1}{2}} \begin{array}{l}
q+r, q_{1}+r \\
p+q_{1}+k, p_{1}+q+k
\end{array}\left[\begin{array}{l}
\left.x^{2 \gamma} \left\lvert\, \begin{array}{l}
-{ }_{1} d_{q_{1}},{ }_{1} \sigma_{k},-{ }_{1} a_{p} \\
{ }_{1} b_{q},-{ }_{1} \sigma_{k},{ }_{1} \sigma_{1} p_{1}
\end{array}\right.\right] \\
q+q_{1}+4 r>p+p_{1}+2 k>2 r
\end{array},\right.
\end{array}
$$

which transforms an $R\left(a_{p} ; b_{q}\right)$ into an $R\left(c_{p_{1}} ; d_{q_{1}}\right)$,

$$
\begin{aligned}
& P_{2}(x)=x^{\gamma-\frac{1}{2}} \begin{array}{l}
q+q_{2}+t, t \\
p+p_{2}+s, q+q_{2}+s
\end{array}\left[\begin{array}{l|l}
\beta^{2} x^{2 \gamma} & \left.\begin{array}{l}
\rho_{s},-{ }_{1} a_{p},-{ }_{1} \gamma_{p_{2}} \\
1_{q} b_{q},{ }_{1} q_{2}, \\
q_{1},-\rho_{s}
\end{array}\right],
\end{array}\right. \\
& 2 t<2 s+p+p_{2}<4 t+q+q_{2},
\end{aligned}
$$

transforming an $R\left(a_{p} ; b_{q}\right)$ into an $R\left(\gamma_{p_{2}} ; \delta_{q_{2}}\right)$ and vice.versa.
Substitutlng these kernels in (3•1), replacing $y^{2 \gamma}$ by $y$ and then evaluating the integral thus obtained with the help of a known result [4, p. 422(14)], we get the resultant kernel

$$
\begin{align*}
& \begin{aligned}
\\
K
\end{aligned} \quad \begin{array}{l}
q+q_{1}+q_{2}+t+r, q+t+r \\
G_{2}(1 / \gamma) x^{\gamma-\frac{1}{2}} \\
\\
p+p_{1}+p_{2}+q+s+k, p+q+q_{1}+q_{2}+s+k
\end{array}
\end{align*}
$$

as a kernel transforming an $R\left(c_{p_{1}} ; d_{q_{1}}\right)$ into an $R\left(\gamma_{p_{2}} ; \delta_{q_{2}}\right)$ and vice-versa, provided that $1 \leqslant q_{1}+r \leqslant p+q_{1}+k<p_{1}+q+k<q+q_{1}+q_{2}+k-p_{2}$, $q+q_{1}+4 r>p+p_{1}+2 k>2 r, 0 \leqslant 2 t<p+p_{2}+2 s<q+q_{2}+4 t$, $\frac{1}{2}\left(p+p_{1}+q-q_{1}\right)-r+k<q+r \leqslant p_{1}+q+\dot{k}, \quad \frac{1}{2}\left(p+p_{2}+q_{2}+q\right)+s-t<$ $q+q_{2}+t \leqslant q+q_{2}+s,\left|\arg \beta^{2} x^{2 \gamma}\right|<\left[\frac{1}{2}\left(q+q_{2}-p-p_{2}\right)+2 t-s\right] \pi$, $\operatorname{Re}\left(m_{j}+n_{h}\right)>-1,\left(j=1, \ldots, q+r ; h=1, \ldots \ldots, q+q_{2}+t\right)$ where $m_{j}=b_{j}(j=1, \ldots, q), m_{q+j}=-\sigma_{j}(j=1, \ldots, r), n_{h}=b_{h}(h=1, \ldots, q)$,

$$
\begin{aligned}
& n_{q+k}=\delta h\left(h=1, \ldots, q_{2}\right), n_{q+q_{2}}+h-\rho_{h}(h=1, \ldots, t) \text { and } \operatorname{Re}\left(r_{j}+s_{h}\right)<1 \\
& \left(j=1, \ldots, g_{1}+r ; h=1, \ldots, t\right) \text { where } r_{j}=-d_{j}\left(j=1, \ldots, q_{1}\right), r_{q_{1}+j}=\sigma_{j} \\
& (j=1, \ldots, r), s_{h}=\rho_{h}(h=1, \ldots, t) .
\end{aligned}
$$

3.2.1. Particular cases: (i) In (3.8), taking $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=1$, $q_{1}=q_{2}=q=2, r=0, t=0, s=0, k=0, b_{1}=v / 2, b_{2}=0, a_{1}=0, c_{1}=k-m-\frac{1}{2}-\mu / 2$, $d_{1}=\mu / 2, d_{2}=2 m+\mu / 2, \gamma_{1}=l-n-\frac{1}{2}-\rho / 2, \delta_{1}=\rho / 2, \delta_{2}=2 n+\rho / 2$ and using the identity [3, p. $209(7)]$ we arrive at a known kernel [10, p. 94], transforming an $R_{\mu}(k, m)$ into an $R_{\rho}(l, n)$ and vice versa.
(ii) In (3.8), putting $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=1, q_{1}=q_{2}=q=2, r=0$, $t=0, s=0, k=0, b_{1}=-1, a_{1}=0, c_{1}=k-m-\frac{1}{2}-\mu / 2, d_{1}=\mu / 2, d_{2}=2 m+\mu / 2$, $\gamma_{1}=l-n-\frac{1}{2}-\rho / 2, \delta_{1}=\rho / 2, \delta_{2}=2 n+\rho / 2$ and applying the identities [3, p. 209(7) and p. 215(2)] we get a kernel transforming an $R_{\mu}(k, m)$ into an $R_{\rho}(l, n)$ and vice-versa.
(iii) With $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=0, q_{1}=q_{2}=q=1, r=0, t=0, s=0$, $k=0, b_{1}=-1, d_{1}=\mu / 2$, and $\delta_{1}=\rho / 2$, (3.8) yields a known kernel [8, p. 59] which transforms an $R_{\mu}$ into an $R_{\rho}$.
(iv) Similarly, setting the parameters suitably in (3.8), it is easy to deduce the kernels due to Narain [8, p. 60] and Saxena [10, p. 94] transforming an $R_{\mu}$ into an $R_{\mu-1}$ and an $R_{\mu}$ into an $R_{\rho}$ respectively.
4. Theorem 2 : The resultant

$$
K(x)=\int_{0}^{\infty} P_{1}(x y) P_{u}(y) d y,
$$

of the two kernels of the form :
and
where $\quad 0<c<1, \quad 0<c^{\prime}<1$ and $\chi(s)=\chi(1-s) ; \quad \omega(s)=\omega(1-s)$, is a kernel transforming an $R\left(c_{p_{1}} ; d_{q_{1}}\right)$ into an $R\left(a_{p_{2}} ; \beta_{q_{2}}\right)$.

Its proof follows, if we proceed on similar lines as those of theorem 1
4.1. Corollarres: (i) Having $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=1, q_{1}=q_{2}=q=2$, $a_{1}=k^{\prime}-m^{\prime}-\frac{t}{2}-\mu / 2, b_{1}=\mu / 2, b_{2}=2 m^{\prime}+\mu / 2, c_{1}=k-m \frac{1}{2}-\nu / 2, d_{1}=\nu / 2, d_{2}=2 m+\nu / 2$, $\alpha_{1}=l-n-\frac{1}{2}-\lambda / 2, \beta_{1}=\lambda / 2$, and $\beta_{2}=2 n+\lambda / 2$, the above theorem reduces to a result, given by Saxena [10, p. 95].
(ii) With $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=0, q_{1}=q_{2}=q=1, b_{1}=\mu / 2, d_{1}=\nu / 2$ and $\beta_{1}=\lambda / 2$, the theorem yields a known result due to Agrawal [1, p. 308].
4.2. Example: Let us have the kernels [11]

transforming an $R\left(a_{p} ; b_{q}\right)$ into an $R\left(c_{p} ; d_{q_{1}}\right)$ and vice-versa ;
and
$P_{2}(x)=x^{\gamma-t} G^{q+q_{3}+t, t}$
$p+p_{2}+s, q+q_{2}+s$$\left[\begin{array}{l|l}\beta^{2} x^{2 \gamma} & \left.\begin{array}{c}{ }^{\rho_{s},-{ }_{1}} a_{p},-{ }_{1} \gamma_{p} \\ { }_{1} b_{q},{ }_{1} \delta_{2},-{ }_{2} \rho_{s}\end{array}\right], 2 t<2 s+p+p_{2}<4 t+q+q_{2},\end{array}\right.$ transforming an $R\left(a_{p} ; b_{q}\right)$ into an $k\left(\gamma_{p_{2}} ; \delta_{q_{2}}\right)$ and vice-versa.

Substituting in ( $4 \cdot 1$ ), replacing $y^{2 \gamma}$ by $y$ and then evaluating the integral thus obtained with the help of the known result [4, p. 422(14)] we get the resultant kernel

$$
\begin{align*}
& K(x)=\frac{x^{\gamma-\frac{1}{2}}}{2 \gamma \beta^{2}} \quad G^{q+q_{1}+r+t, q+q_{3}+r+t} \quad p+p_{1}+q+q_{2}+k+s, p+p_{2}+q+q_{1}+k+s
\end{align*}
$$

transforming an $R\left(c_{p_{1}} ; d_{q_{1}}\right)$ into an $R\left(y_{p_{2}} ; \delta q_{q_{2}}\right)$, provided that $0 \leqslant t \leqslant p+p_{2}+$ $s<q+q_{2}+s<p_{2}-p_{1}+q+q_{1}+s, 0 \leqslant 2 r<2 k+p+p_{1}<4 r+q_{1}+q, 2 t<2 s+p+p_{2}$ $<4 t+q+q_{1}, \frac{3}{2}\left(p+p_{2}+q+q_{2}\right)+s-t<q+q_{2}+t \leqslant q+q_{2}+s, \frac{1}{2}\left(p+p_{1}\right.$ $\left.+q+q_{1}\right)+k-r<q+q_{1}+r \leqslant q+q_{1}+k,\left|\arg \beta^{2} x^{2 \gamma}\right|<\left[\frac{1}{2}\left(q+q_{1}-p-p_{1}\right)\right.$ $+2 r-k] \pi ; R\left(m_{j}+n_{h}\right)>-1\left(j=1, \ldots, q+q_{2}+t ; h=1, \ldots, q+q_{1}+r\right)$ where $m_{j}=b_{j}(j=1, \ldots, q), m_{q+j}=\delta_{j}\left(j=1, \ldots, q_{2}\right), m_{q+q+j}=-\rho_{j}(j=1, \ldots t)$, $n_{h}=b_{h}(j=1, \ldots, q), n_{q+h}=d_{h}\left(h=1, \ldots, q_{1}\right), n_{q+q_{1}+h}=-\sigma_{h}(h=1, \ldots, 1) ;$ $\operatorname{Re}\left(\rho_{j}+\sigma_{h}\right)<1(j=1, \ldots, t ; h=1, \ldots, r)$
and $\sum_{1}^{s} \operatorname{Re} \rho_{j}-\sum_{1}^{p} \operatorname{Re} a_{j}-\sum_{1}^{p_{2}} \operatorname{Re} \gamma_{j}-\sum_{1}^{q} \operatorname{Re} b_{j}-\sum_{1}^{q_{2}} \operatorname{Re} \delta_{j}+\sum_{1}^{s} \operatorname{Re} \rho_{j}+\frac{1}{2}\left(q+q_{2}-p-p_{i}+1\right)$ $>\left(q+q_{2}-p-p_{2}\right) \cdot \operatorname{Re}\left(\sigma_{j}\right)(j=1, \ldots, r)$.
$4 \cdot 2 \cdot 1$. Particular cases: (i) In (4.2), setting $\beta=\downarrow, \gamma=1, p_{1}=p_{2}=p=1$, $q_{1}=q_{2}=q=2, r=0, t=0, s=0, k=0, b_{1}>\mu_{1} 2, b_{2}=0, a_{1}=0, c_{1}=k-m-\frac{1}{2}-\mathrm{r} / 2$,
$d_{1}=v / 2, d_{2}=2 m+\nu / 2, \gamma_{1}=l-n-\frac{1}{2}-\lambda / 2, \delta_{1}=\lambda / 2, \delta_{2}=2 n+\lambda / 2$ and using a known identity [3, p. 209(7)], we get a known kernel (10, p. 98] transforming an $R_{\nu}(k, m)$ into an $R_{\lambda}(l, n)$.
(ii) In (4•2), putting $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=0, q_{1}=q_{2}=q=1, t=1$, $s=2, r=0, k=0, a_{1}=0, c_{1}=0, \gamma_{1}=0, b_{1}=\mu / 2, d_{1}=1+\nu / 2, \delta_{1}=v / 2$, $\rho_{1}=-\nu / 2, \rho_{2}=1+\nu / 2$, we obtain a kernel [1, p. 309] transforming an $R_{\nu+2}$ to an $R_{\nu}$.
(iii) In, (4.2), having $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=0, q_{1}=q_{2}=q=1$, $r=0, t=0, s=0, k=0, a_{1}=0, c_{1}=0, \gamma_{1}=0$ and either $\left(b_{1}=\mu / 2, \delta_{1}=v / 2\right.$, $d_{1}=\nu / 2$ ) or ( $b_{1}=\nu / 2, \delta_{1}=\mu / 2, d_{1}=\mu / 2$ ), we arrive at the kernel [1, p. 310] transforming an $R_{\nu}$ or an $R_{\mu}$ into itself respectively.
5. Theorem 3: The resultant

$$
K(x)=\int_{0}^{\infty} P_{1}(y) p_{2}(x y) d y,
$$

of the two kernels of the form :
and
where $0<c<1,0<c^{\prime}<1$ and $\chi(s)=\chi(1-s), \omega(s)=\omega(1-s)$, is a kernel transforming an $R\left(c_{p_{1}} ; d_{q_{1}}\right)$ into an $R\left(\alpha_{p_{2}} ; \beta q_{2}\right)$.

The proof of this theorem can be developed on the similar lines as those of the theorem 1 .

5•1. Corollaries: (i) Setting $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=1, q_{1}=q_{2}=q=2$, $a_{1}=k^{\prime}-m^{\prime}-\frac{1}{2}-\nu^{\prime} / 2, b_{1}=\nu^{\prime} / 2, b_{2}=2 m^{\prime}+\nu^{\prime} / 2, c_{1}=l-n-\frac{1}{2}-\mu^{\prime} / 2, d_{1}=\mu^{\prime} / 2$, $d_{2}=2 n+\mu^{\prime} / 2, \alpha_{1}=k-m-\frac{1}{2}-\rho / 2, \quad \beta_{1}=\rho / 2$, and $\beta_{2}=2 m+\rho / 2$, the above theorem reduces to a result, given by Saxena [10, p. 99].
(ii) With $\beta=\frac{1}{2}, \quad \gamma=1, p_{1}=p_{2}=p=0, \quad q_{1}=q_{2}=q=1, \quad b_{1}=\nu^{\prime} / 2$, $d_{1}=\mu^{\prime} / 2$ and $\beta_{1}=\rho / 2$, the theorem yields a known result due to Agarwal [1, p.309].
52. Example: Let us choose the kernels [11].
$P_{1}(x)=x^{\gamma-\frac{1}{2}} \begin{aligned} & q+r_{2} q_{1}+r \\ & p+q_{1}+k, p_{1}+q+k\end{aligned}\left[x^{a \gamma} \left\lvert\, \begin{array}{l}-{ }_{1} d_{q_{1}},{ }_{1} \sigma_{k} \sigma_{2}-{ }_{1} a_{p} \\ 1_{q},{ }_{1} \sigma_{k},{ }_{1} c_{p}\end{array}\right.\right], q+q_{1}+4 r>p+p_{1}+2 k>2 r$, transforming an $R\left(a_{p} ; b_{q}\right)$ into an $R\left(c_{p_{1}} ; d_{q_{1}}\right)$;
and
 transforming an $R\left(a_{p} ; b_{q}\right)$ into an $R\left(\gamma_{p_{2}} ; \delta q_{2}\right)$.

Putting these in ( $5 \cdot 1$ ), replacing $y^{2 \gamma}$ by $y$ and then evaluating the integral thus obtained with the help of the known result [4, p. 422(14)] we obtain the resultant kernel

$$
\begin{align*}
& K(x)=\begin{aligned}
& x^{\gamma-1} \\
& 2 \gamma G^{q+q_{1}+r+t, q+q_{2}+r+t} \\
& p+p_{1}+q+q_{2}+k+s, p+p_{2}+q+q_{1}+k+s
\end{aligned} \\
& {\left[\begin{array}{l|l}
x^{a \gamma} & -{ }_{1} b_{q},{ }_{1} \sigma_{r},-{ }_{1} \delta_{q_{2}},{ }_{1} \rho_{s},-{ }_{1} a_{p}, r+{ }^{\sigma}{ }_{k},-{ }_{1} c_{p} \\
{ }_{1} d_{q_{1}},-{ }_{1} \sigma_{r},{ }_{1} b_{q},-{ }_{1} \rho_{s, 1} \gamma_{p}, & -r+{ }_{2} \sigma_{k},{ }_{1} a_{p}
\end{array}\right],}
\end{align*}
$$

as transforming an $R\left(c_{p_{1}} ; d_{q_{1}}\right)$ into an $R\left(\gamma_{p_{2}} ; \delta q_{2}\right)$, provided that

$$
\begin{aligned}
& 1 \leqslant q_{1}+r \leqslant p+q_{1}+k<p_{1}+q+k<q+q_{1}-q_{2}+p_{2}+k, 0 \leqslant q_{2}+t \leqslant p \\
& +q_{2}+s, 2 r<p+p_{1}+2 k<q+q_{1}+4 r, 2 t<p+p_{2}+2 s<q+q_{2}+4 t, \\
& \frac{1}{2}\left(p+p_{1}+q-q_{1}\right)+k-r<q+r \leqslant p_{1}+q+k, \frac{1}{2}\left(p+p_{2}+q q_{2}\right)+s-t \\
& <q+t \leqslant p_{2}+q+s,\left|\arg x^{2} \lambda\right|<\left[\frac{1}{2}\left(q+q_{2}-p-p_{2}\right)+2 t-s\right] \pi ; \operatorname{Re}\left(m_{j}+n_{h}\right) \\
& >-1(j=1, \ldots, q+r ; h=1, \ldots, q+t) \text { where } m_{j}=b_{j}\left(j=1, \ldots,{ }_{2}\right), \\
& m_{q+j}=-\sigma_{j}(j=1, \ldots, r), n_{h}=b_{h}(h=1, \ldots, q), n_{q+h}=-\rho_{h}(h=1, \ldots, t) ; \\
& \operatorname{Re}\left(r_{j}+s_{h}\right)<1\left(j=1, \ldots \ldots, q_{1}+r ; h=1, \ldots, q_{2}+t\right) \text { where } r_{j}=-d_{j} \\
& \left(j=1, \ldots, q_{1}\right) r_{q_{1}+j}=\sigma_{j}(j=1, \ldots, r), s_{h}=-\delta h\left(h=1, \ldots, q_{2}\right), s_{q_{2}+h}=o_{h} \\
& (h=1, \ldots, t) .
\end{aligned}
$$

5.2.1. Particular cases: (i) In (5.2), having $\beta=\frac{1}{2}, \gamma=1, p_{1}=p_{2}=p=0$, $q_{1}=q_{2}=q=1, r=1, k=2, t=0, s=0, a_{1}=0, c_{1}=0, \gamma_{1}=0, b_{1}=\mu / 2$, $d_{1}=1+v / 2, \delta_{1}=v / 2, \sigma_{1}=-v / 2$, and $\sigma_{2}=1+1 / 2$, we again get the kernel given by Agrawal [1, p. 309], transforming an $R_{\nu+2}$ into an $R_{\nu}$.
(ii) Similarly, by setting the parameters suitably, the kernel due to Agrawal [1, p. 310], transforming an $R_{\nu}$ or an $R_{\mu}$ into itself, can easily be deduced as particular cases of (5•2).

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# Application of generalized function method to the theory water waves 

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#### Abstract

The generalized function method is utilized in solving an initial value problem of water waves generated in an inviscid, incompressible and homogeneous ocean by a harmonically oscillating pressure distribution acting on its free surface. It is substantiated that this new approach to the theory of water waves is very simple, rigorous and useful from the point of view of sufficient generality and applicability. An asymptotic analysis is carried out for finding both the steady state and transient solutions of the problem. A discussion of the waves motions is made with a special reference to their charecteristic features and the effects of surface tension on the wave phenomena.


## 1. Introduction

Water wave phenomena have received a great deal of attention by Stoker (1), Miles (2), Debnath (3-5), Debnath and Rosenblat (6), Lamb (7-8), Sen (9) and many others. Several classical methods have so fir been employed in sludying the linearized theory of wave motions in an inviscid, incompressible homogeneous fluid with a free surface due to local disturbances acting on and below the free surface of the fluid. An extensive use of classical Fourier series and transform techniques in these problems are well known and readily available in the existing literature (Wehausen and Laitone (10), Lamb (8), and Stoker (1)) on the subject.

In recent years, the authur (3-5) has clearly pointed out certain inherent difficulties involved in the steady state water wave problems and the classical Fourier analysis used in solving them. These difficulties are essentially related to the existence, uniqueness and the real singularities of the inversion Fourier integral representation of the solution. In order to resolve these difficulties, Debnath (3) alone and in collaboration with Rosenblat (6) have investigated the steady state and the transient problem of water wave phenomena due to a harmonically oscillating pressure distribution of quite general character. This new treatment enables them to derive a mathematically unique and physically realistic solution in a simple, straightforward and elegant manner.

It thus appears that the water wave problems in various situations are readily and rigorously accessible to the generalized function treatment. This approach can not only eliminate the deficiencies of the classical Fourier transform analysis, but also give the solution of physical interest. Thus it seems to be extremely useful device from the point of view of sufficient generality and applicability. Another convincing point about this method is that there is neither any need for modification of the basic flow equations nor any justification of the limit operation suggested by Lamb and Sen.

This paper is devoted to study of water waves produced in an inviscid, incompressible, homogeneous deep ocean by a harmonically oscillating pressure distribution of general nature "switched on" at time $t=0+$. The aim of this work is two fold. In the first place, the study provides an interesting and useful application of the generalized functions to the water wave problems. Secondly, a unique solution of physical interest is obtained rigorously without any help of a radiation condition or an equivalent device. An asymptotic analysis is carried out to investigate the steady state and the transient solutions. A discussion of the wave motions is made with a special reference to their characteristic features and the effects of surface tension on the wave phenomena. Some consequencesland implications of the results are presented.

## 2. Statement of the initial value problem

We consider a linearized two-dimensional transient wave motion of an inviscid, incompressible, homogeneous fluid of infinite depth. We take the origin of the rectangular cartesian coordinates $(x, y, z)$ on the free surface of the fluid of uniform density $\rho$ and surface tension $T$. In the undisturbed state, $x-y$ plane is assumed to be horizontal free surface and $z$-axis vertical positive upwards.

The unsteady wave motion is. set up on the free surface of the fluid by an oscillating periodic pressure of fixed frequency $\omega$ in the form

$$
\left.\begin{array}{rlrl}
p(x, t) & =P p(x) e^{i \omega t} & H(t),-a \leqslant x \leqslant a \\
& =0 & , \text { outside }
\end{array}\right\}
$$

on the undisturbed free surface $z=0$ of the fluid, where $P$ is a constant, $p(x)$ is any arbitrary function of $x$ and $H(t)$ is the Heaviside unit step function of time $t$.

In view of the fact that the flow is irrotational, there exists a wave potential $\phi(x, z ; t)$ which is governed by the Laplace equation

$$
\nabla^{2} \phi \equiv \phi_{x: x}+\phi_{z z},-\infty<z \leqslant 0,-\infty<x<\infty,
$$

where suffixes stand for the partial differentiation.
The linearized dynamic and kinematic free surface conditions are given by

$$
\begin{gather*}
\phi_{t}+g \eta-\frac{T}{\rho} \eta_{x x}=-\frac{P}{\rho} p(x) e^{i \omega t}, z=0, t>0, \\
\eta_{t}=\phi_{z}, z=0, t>0,
\end{gather*}
$$

where $\eta=\eta(x, t)$ denotes the vertical surface elevation and $g$ is the gravitational acceleration.

The bottom boundary condition is given by

$$
\phi_{2} \rightarrow 0 \text { as } z \rightarrow-\infty, t>0,
$$

The initial conditions are

$$
\left.\begin{array}{r}
\phi(x, z, 0)=\eta(x, 0)=0 \\
{\left[\phi_{l}(x, 0 ; t)\right]_{t=0}=-\frac{1}{\rho} p(x, 0)}
\end{array}\right\}
$$

In addition, we shall treat $\phi(x, z ; t)$ and $\eta(x, t)$ as generalized functions of $x$ in the sense of Lighthill (11) so that their Fourier transform exists with respect to $x$.

## 3. Integral solutions

It is convenient to use the joint Laplace and the generalized Fourier trans* formations defined by the integral

$$
\vec{\phi}=\widetilde{\bar{\phi}}(k, z ; s)=\frac{1}{\sqrt{2 \pi}} \int_{-\infty}^{\infty} e^{-i k x} d k \int_{0}^{\infty} e^{-s t} \phi(x, z ; t) d t
$$

The joint transform enables us to simplify the equations (2.2)-(2.5) and it turns out that the solution for the transform functions is

$$
\begin{align*}
& \overline{\tilde{\phi}}(k, z ; s)=-\frac{P \bar{p}(k) s e^{|k| z}}{\rho} \frac{(s-i \omega)\left(s^{2}+a^{2}\right)}{} \\
& \bar{\sim} \\
& \left.\eta(k, s)=-\frac{p}{\rho}(s-i \omega) \right\rvert\, \frac{k \mid}{\left(s^{2}+a^{2}\right)}
\end{align*}
$$

where $\quad \alpha^{2} \equiv \alpha^{2}(k)=\left(g+\frac{T}{\rho} k^{9}\right)|k|$,
The inversion theorem of the Laplace and the Fourier transform together with the Fultung theorem of the former yields

$$
\begin{gather*}
\phi(x, z ; t)=\frac{P}{\rho \sqrt{2 \cdot \pi}} \int_{-\infty}^{\infty}\left(i \omega \cos \alpha t-\alpha \sin \alpha t-i \omega e^{i \omega t}\right) \frac{\bar{p}(k) e^{|k| z+i k x}}{\left(x^{2}-\omega^{2}\right)} d k,(35) \\
\eta(x, t)=\frac{P}{\rho \sqrt{2 \pi}} \int_{-\infty}^{\infty}|k| \bar{p}(k)\left(\alpha \cos \alpha t-i \omega \sin \sim t-\alpha e^{i \omega t}\right) \frac{e^{i k x}}{\alpha\left(\alpha^{2}-\omega^{2}\right)} d k,
\end{gather*}
$$

Usually these integrals can not be evaluated exactly and it is necessary to resort to asymptotic techniques.

## 4. Asymptotic analysis

It would be sufficient for the examination of the transient wave motions to evaluate either the integral $(3 \cdot 5)$ or $(3 \cdot 6)$. We next turn our attention to the evaluation ( $3 \cdot 6$ ) by using formulae for the asymptotic development of the generalized Fourier transform due to Lighthill (11) and Jones (12). In order to facilitate the analysis, it is convenient to set $T=0$ in (3.6) and to include a discussion about the effects of surface tension on the wave motions in a subsequent section.

Integral (3.6) can be written as

$$
\eta(x, t)=\frac{P}{2 g \rho \sqrt{2 \pi}}(T-S)
$$

where $S$ and $T$ represent the steady state and the transient wave integrals respectively and are given by

$$
\begin{align*}
& S=e^{i \omega t} \int_{-\infty}^{\infty} \alpha(k) \bar{p}(k)\left(\frac{1}{i \alpha+\omega}+\frac{1}{a-\omega}\right) e^{i k x} d k, \\
& T=\int_{-\infty}^{\infty} \alpha(k) \bar{p}(k)\left(\frac{e^{-i \alpha t}}{\alpha-\omega}+\frac{e^{i a t}}{\alpha-\omega}\right) e^{i k x} d k,
\end{align*}
$$

where $\alpha$ is given by $(3 \cdot 4)$ with $T=0$.
It may be noted that the significant contributions to the steady state integral $S$, as $|x| \rightarrow \infty$, come from the poles of the integrand of (4.2), that is, from the points where $a \pm \omega=0$. In fact, the poles are at $k= \pm \frac{\boldsymbol{\omega}^{2}}{g}$. We evaluate $S$ by using the result (11, pages 43 and 52) and it turns out that as $|x| \rightarrow \infty$

$$
S \sim \pi i \operatorname{sgn} x \psi\left(\frac{\omega^{2}}{S}\right)\left(e^{i \frac{\omega^{2} x}{g}}-e^{-i \frac{\omega^{2} x}{g}}\right)+0\left(\frac{1}{\mid x}\right)
$$

where

$$
\psi\left(\frac{\omega^{2}}{g}\right)_{--}=\frac{\alpha\left(\frac{\omega^{2}}{g}\right) \bar{p}\left(\frac{\omega^{2}}{g}\right)}{\left(\frac{d a}{d k}\right)_{k=\omega^{2} / g}}
$$

with $\bar{p}(k)=\bar{p}(-k)$ which corresponds to some simple cases of physical interest. However, this even-character of $\bar{p}^{-}(k)$ is not necessary in a general situation.

The method of stationary phase (Jone, 12) in conjunction with the result stated above can be used for the asymptotic evaluation of (4.3) for large $t$. It is then necessary to locate the stationary points of (4.3) which are roots of equation $\frac{d \alpha}{d k}= \pm \frac{x}{t}$ and $k= \pm k_{1}$ with $k_{1}=\frac{g t_{r}^{2}}{4 x^{2}}$. Thus the transient contribution to (4.3) is obtained in the form

$$
\begin{align*}
& T \sim a\left(k_{1}\right) \bar{p}\left(k_{1}\right)\left\{\frac{2 \pi}{t\left|a^{\prime \prime}\left(k_{1}\right)\right|}\right\}^{\frac{1}{2}}\left[\frac{e-i\left\{a\left(k_{1}\right) t+k_{1} x+\frac{\pi}{4}\right\}}{a\left(k_{1}\right)+\omega}\right. \\
&\left.=\frac{e^{i\left\{t \alpha\left(k_{1}\right)-k_{1} x-\frac{\pi}{4}\right\}}}{a\left(\overline{\left.k_{1}\right)-\omega}\right.}\right]+0\left(\frac{1}{t^{2}}\right)
\end{align*}
$$

It remains to calculate the contribution to $T$ from its polar singularities which are the same as those of $S$. To evaluate the polar contribution to $T$, we write

$$
T_{\text {polar }}=T_{1}+T_{2}
$$

where $T_{1}$ and $T_{2}$ are given by

$$
\begin{align*}
& T_{1}=\int_{-\infty}^{\infty} \frac{a(k) \bar{p}(k) e^{-i(\alpha t-k x) d k}}{\alpha+\omega}  \tag{4}\\
& T_{2}=\int_{-\infty}^{\infty} \frac{\alpha(k) \bar{p}(k) e^{i(a t+k x)} d k}{\alpha-\omega}
\end{align*}
$$

It may be noted that $T_{1}$ has no poles in the range of integration and hence contributes nothing. On the other hand, $T_{2}$ has two poles at $k= \pm t \frac{\omega^{2}}{g}$ and the residue-contribution can easily be obtained. Making reference to the technique developed by Debnath and Rosenblat (6) together with a change of variable $k$ to $a(k)$, it turns out that

$$
T_{\text {polar }}=\pi i e^{i \omega t} \psi\left(\frac{\omega^{2}}{g}\right)\left(e^{i \frac{\omega^{2}}{g} x}+e^{-i \frac{\omega^{2}}{g} x}\right),
$$

for large $t>0$. The multiplicative factor sgnt is omitted since $t>0$ is the interest.

The asymptotic solution for the veritical surface elevation $\eta(x, t)$ is obtained in the form

$$
\begin{gather*}
\eta(x, i) \sim \frac{2 P i \omega^{2} \pi \bar{p}\left(\frac{\omega^{2}}{g}\right) i\left(\omega t-\frac{\omega^{2}}{g}|x|\right)}{\rho g^{2} \sqrt{2 \pi}} e^{e} \begin{array}{l}
P t^{2} \bar{p}\left(\frac{g t^{2}}{4 x^{2}}\right) \sqrt{g}-\left[\frac{g t}{2 x \omega} \cos \left(\frac{g t^{2}}{4 x}-\frac{\pi}{4}\right)+i g \sin \left(\frac{g t^{2}}{4 x}-\frac{\pi}{4}\right)\right] \\
+\frac{2 \sqrt{2} \rho \omega x^{5 / 2}\left[\binom{g t}{2 x \omega}^{2}-1\right]}{\sqrt{2}]}+0\left(\frac{1}{t^{2}}\right),
\end{array},
\end{gather*}
$$

Remark: Solution ( $4 \cdot 11$ ) breaks down at $g t=2 x \infty$ which may be regarded as the critical point. A special device advanced by Debnath (13) may be employed to find a solution valid at the critical point. As the solution for large $t$ is of most interest, no attention is given to this point.

## 5. Some pressure distribution of physical interest

The entire analysis has been carried out for an arbitrary function $p(x)$ involved in (2•1). It would be sufficient for the investigation of the dominant features of the unsteady wave motions to take some simple form of $p(x)$, such as,
(a) $p(x)=\left(a^{2}-x^{2}\right)^{-\frac{1}{2}}$,
(b) $p(x)=\frac{c^{2}}{c^{2}+x^{2}}$,
(c) $p(x)=1$,
(d) $p(x)=d^{-x^{2} / x_{0}^{2}}$ and
(e) $p(x)=\delta(x)$,
where $\delta(x)$ is the Dirac function of distribution.
The corresponding Fourier transform $\bar{p}(k)$ of these cases $(a)-(e)$ are respectively given by
(a) $\bar{p}(k)=\sqrt{\frac{\pi}{2}} J_{0}(a k)$,
(b) $\bar{j}(k)=\sqrt{\frac{2}{\pi}} c e^{-|k| c}$,
(c) $\sqrt{\frac{2}{\pi}} \frac{\sin a k}{k}$
(d) $e^{-k^{2} m^{2}},\left(m^{2}=\frac{x_{0}{ }^{2}}{4}\right)$ and
(e) $\frac{1}{\sqrt{2 \pi}}$

$$
+\frac{P t^{2} \exp \left\{-m^{2}\left(\frac{g t^{2}}{4 x^{2}}\right)^{2}\right\}\left[\frac{g t}{2 x \omega} \cos \left(\frac{g t^{2}}{4 x}-\frac{\pi}{4}\right)+i \sin \frac{g t^{2}}{4 x}-\frac{\pi}{4}\right]}{2 \sqrt{2} \rho \omega x^{5 / 2}\left[\left(\frac{g t^{2}}{2 x \omega}\right)^{2}-1\right]}
$$

This solution clearly suggests that its transient component decays to zero very rapidly as $t \rightarrow \infty$ for fixed $x$. Consequently, an ultimate steady state is established in the limit. An exactly same conclusion can be drawn in all other pressure distributions confined over a finite region of the free surface. The nature of ( $5 \cdot 1$ ) in non-dimensional form at fixed time $\iota^{*}=25 \pi$ is exhibited by Figure 1 . The non-dimensional quantities $\eta^{*}, x^{*}, t^{*}$ are defined by $\eta^{*}=\frac{P \omega^{4}}{\rho g^{3}} \eta, t^{*}=\omega t$ and $x^{*}=\frac{\infty^{2}}{g} x$.

The asymptotic solution for $\eta(x, t)$ related to case (e) has the representation

$$
\begin{align*}
& \eta(x, t) \sim \frac{P i \omega^{2}}{\rho g^{2}} e^{i\left(\omega t-\frac{\omega^{2}}{g}|x|\right)} \\
& +\frac{P \sqrt{g} t^{2}\left[\frac{g t}{2 x \omega} \cos \left(\frac{g t^{2}}{4 x}-\frac{\pi}{4}\right)+i \sin \left(\frac{g t^{2}}{4 x}-\frac{\pi}{4}\right)\right]}{4 \sqrt{\pi^{-} \rho \omega} x^{\mathrm{b} / 2}\left[\left(\frac{g t}{2 x \omega}\right)^{2}-1\right]}+0\left(\frac{1}{t^{2}}\right)
\end{align*}
$$

It may be observed that this result follows fiom that of case (c) as a limit $P \rightarrow \infty a \rightarrow 0$ provided $2 P a$ tends or a finite constant.


## 6. Conclusions

It may be inferred from the above asymptotic solution that the ultimate steady state is attained in the limit $t \rightarrow \infty$ for fixed $x$. The steady state solution for the surface elevation $\eta(x, t)$ has the form

$$
\eta(x, t) \sim \frac{\sqrt{2 \pi} P i \omega^{2}}{\rho g^{2}} \bar{p}\left(-\omega^{2}\right) e^{i\left(\omega t-\frac{\omega^{2}}{g}|x|\right)}
$$

This solution represents progressive surface waves propagating with the phase velocity $\frac{g}{\omega}$ and the group velocity $\frac{g}{2 \omega}$, and the amplitude of the waves is constant.

This analysis shows that the present method of solution provides an interesting example of the application of generalized function method in the theory of water waves. Furthermore, a comparison of this method with the classical methods suggests that the use of generalized functions simplifies the analysis considerably.

## 7. Effects of surface tension

It has already been seen that surface tension enters into the water wave phenomena through equation (2.3). Consequently, the polar singularities of the wave integrals are at the points where

$$
s^{2}+\left(g+\frac{T}{\rho} k^{2}\right)|k|=0
$$

Naturally, the stationary points related to the transient integrals suffer from a change due to inclusion of surface tension.

Without entering into mathematical details, it may be noted that real roots of $(7 \cdot 1)$ and the stationary points associated with the transient integrals are given by

$$
\begin{align*}
& k=k_{0}^{*}=\frac{\omega^{2}}{g}-\varepsilon, \varepsilon=\frac{T}{\rho}\left(\frac{\omega^{3}}{g^{2}}\right)^{2}>0 \\
& k=k_{1}^{*}=\frac{g t^{2}}{4 x^{2}}\left(1+\frac{5 T}{\rho}\right)
\end{align*}
$$

where surface tension $T$ is very small.
Some conclusion can now be drawn concerning the role of surface tension. In the first place, both the steady state and the transient solutions are independently modified by the surface tension. The phase and the group velocities are $\frac{\omega}{k_{0}{ }^{*}}$ and $\frac{\Phi}{2 k_{0}{ }^{*}}$ respectively and evidently are greater than those of the corresponding results obtained without including the effects of surface tension. Secondly, it is of special interest to point out that the transient solution related to case $5(e)$ is $0\left(t^{-4}\right)$ and hence tends to zero as $t \rightarrow \infty$ for fixed $x$. Consequently, the ultimate steady state is reached in the limit. Like viscosity, surface tension is really responsible for the attainment of the steady state.

Physically, the principal effects of surface tension is to increase the velocity of energy propagation and makes the energy more readily distributed among the rapidly travelling waves. It is thus highly likely that the intense accumulations of energy are no longer possible in the far field as $t \rightarrow \infty$. This is one of the most significant contrasts with the corresponding problem which takes on account of surface tension.

## 8. Concluding remarks

It appears from the above analysis that wave phenomena in fluids and elastic solids in various situations can readily and rigorously be investigated with the aid of the generalized function treatment. It may be interesting to analyze the above problem in a more general geometric configuration. Such study will be considered in a subsequent paper.

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## Kinetics and Mechanism of Reaction between Bromine and Quinoline in Carbon-Tetrachloride Solution.

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## Abstract

The product $\left(\mathrm{C}_{9} \mathrm{H}_{7} \mathrm{NBr}_{2}\right)$ of the reaction between bromine and quinoline has been isolated and analysed for available bromine.

Calculated for: $\mathrm{C}_{9} \mathrm{H}_{7} \mathrm{NBr}_{2}$ : $\mathrm{Br}, 55 \cdot 30 \%$
Found : $\quad \mathrm{Br}, 55.00 \%$
The rate equation has been derived as

$$
-\frac{\mathrm{d}\left[\mathrm{Br}_{2}\right]}{\mathrm{dt}}=\frac{2 \mathrm{k}_{1} \mathrm{k}_{2}[\mathrm{Qu}]^{2}\left[\mathrm{Br}_{2}\right]}{\mathrm{k}_{-1}+\mathrm{k}_{2}[\mathrm{Qu}]}
$$

Where $k_{-1}, k_{1}$ and $k_{2}$ are the velocity constant of the different steps of the reaction.

## Introduction

In this paper we have studied the nature of the reaction mechanism and the product formed, when bromine and quinoline react in carbontetra-chloride solution. The experimental part has been published in the Indian J. of Chemistry in 1968, Vol. 6, No. 6, pp. 306-308.

The exact nature of the reaction product when bromine reacts with quinoline in carbontetra-chloride, at ordinary temperature has been discussed by Eisch ${ }^{1,2}$ and the structure formula which has been ascribed to this addition product is as follows :

[ 440 ]

The reaction between bromine and quinoline is very fast and when the concentration of quinoline is comparable to that of bromine (i.e. for initial [bromine] $=15.52 \times 10^{-3} \mathrm{M}$, [quinoline] $=67.58 \times 10^{-3} \mathrm{M}, 33.78 \times 10^{-3} \mathrm{M}$ and $16^{6} 89 \times 10^{-3} \mathrm{M}$ respectively) the total order of the reaction is three -two with respect to quinoline and one with respect to bromine. When the reaction is studied in the presence of excess of quinoline, (i.e. for initial [bromine] $=15 \cdot 52 \times 10^{-3} \mathrm{M}$, [quinoline] $=8.41 \times 10^{-2} \mathrm{M}, 12.62 \times 10^{-2} \mathrm{M}, 16.83 \times 10^{-2} \mathrm{M}$ and $50.50 \times 10^{-2} \mathrm{M}$ respectively and for initial [quinoline] $=42.08 \times 10^{-2} \mathrm{M}$, [bromine] $=27.16$ $\times 10^{-3} \mathrm{M}, 23.28 \times 10^{-3} \mathrm{M}, 13.58 \times 10^{-3} \mathrm{M}$ and $11.64 \times 10^{-3} \mathrm{M}$ respectively) the reaction becomes very fast and also, rapid precipitation of quinoline di-bromide occurs. Under such conditions the reaction becomes second order-one with respect to each of the reactants.

## Study of the Reaction Product

By equivalent method it is found that number of moles of bromine required by one mole of quinoline is one. Also the orange precipitate of the reaction product was collected, washed with fresh carbon-tetrachloride and dried in a desiccator over a mixture of KOH pellets and paraffin shavings. After the orange solid was dried in the dark, it was analysed for "available bromine". The prepared samples were treated with 15 ml . of $40 \%$ KI Solution and 15 ml . of 5 N acetic acid and the liberated iodine was titrated with standard sodium thiosulphate solution.

Calculated for: $\quad \mathrm{C}_{9} \mathrm{H}_{7} \mathrm{NBr}_{2}: \mathrm{Br}, \quad 55.30 \%$
Found : $\quad \mathrm{Br} ; 55.0 \%$
The orange adduct melted at $81-82^{\circ} \mathrm{C}$. Treatment with sodium sulphite and ammonia solution regenerated quinoline. Upon standing in air, it slowly lost bromine vapour and developed yellow flakes.

## Mechanism

We propose the following steps during the reaction between bromine and quinoline in carbon-tetrachloride solution.

$$
\begin{gather*}
\mathrm{Qu} .+\mathrm{Br}_{2} \underset{\mathrm{k}_{-1}}{\stackrel{\mathrm{k}_{1}}{\rightleftharpoons}} \mathrm{Qu} \ldots \mathrm{Br}_{2}{ }^{*}  \tag{1}\\
\mathrm{Qu} \ldots \mathrm{Br}_{2}^{*}+\mathrm{Qu} . \xrightarrow{\mathrm{k}_{2}} 2 \mathrm{Qu} . \mathrm{Br}  \tag{2}\\
\mathrm{Qu} . \mathrm{Br}+\mathrm{Br}_{2} \xrightarrow{\mathrm{k}_{3}} \mathrm{Qu} . \mathrm{Br}_{2}+\mathrm{Br}  \tag{3}\\
\mathrm{Br}+\mathrm{Br} \underset{\mathrm{Br}}{\rightleftharpoons} \mathrm{Br}_{2} \tag{4}
\end{gather*}
$$

The compound $\mathrm{Qr} \ldots \mathrm{Br}_{2}{ }^{*}$ may be regarded as activated complex in which the atoms have not, yet, been organized into a normal molecule, obviously the rate of variation of concentration is given by the expression.

$$
\begin{equation*}
-\frac{\mathrm{d}\left[\mathrm{Br}_{2}\right]}{\mathrm{dt}}=\mathrm{k}_{1}[\mathrm{Qu}]\left[\mathrm{Br}_{2}\right]+\mathrm{k}_{3}\left[\mathrm{Br}_{2}\right][\mathrm{Qu} \mathrm{Br}]-\mathrm{k}_{-1} \quad\left[\mathrm{Qu} \ldots \mathrm{Br}_{2}^{*}\right] \tag{5}
\end{equation*}
$$

From principle of stationary state we have :

$$
-\frac{\left[\mathrm{Qr} \ldots \mathrm{Br}_{2}^{*}\right]}{\mathrm{dt}}=\mathrm{k}_{1}[\mathrm{Qu}]\left[\mathrm{Br}_{2}\right]-\mathrm{k}_{-1}\left[\mathrm{Qu} \ldots \mathrm{Br}_{2}^{*}\right]-\mathrm{k}_{2}\left[\mathrm{Qu} \ldots \mathrm{Br}_{2}^{*}\right][\mathrm{Qu}]=0
$$

or

$$
\begin{equation*}
\left[\mathrm{Qu} \ldots \mathrm{Br}_{2}{ }^{*}\right]=\frac{\mathrm{k}_{1}[\mathrm{Qu}]\left[\mathrm{Br}_{2}\right]}{\mathrm{k}_{-1}+\mathrm{k}_{2}[\mathrm{Qu}]} \tag{6}
\end{equation*}
$$

and

$$
\frac{\mathrm{d}[\mathrm{Qu} \mathrm{Br}]}{\mathrm{dt}}=\mathrm{k}_{2}\left[\mathrm{Qu} \ldots \mathrm{Br}_{2}^{*}\right][\mathrm{Qu}]-\mathrm{k}_{3}[\mathrm{Qu} \mathrm{Br}]\left[\mathrm{Br}_{2}\right]=0
$$

$$
\begin{equation*}
\text { or } \quad[\mathrm{Qu} \mathrm{Br}]=\frac{\mathrm{k}_{2}\left[\mathrm{Qu} \ldots \mathrm{Br}_{2}{ }^{*}\right][\mathrm{Qu}]}{\mathrm{k}_{3}\left[\mathrm{Br}_{2}\right]} \tag{7}
\end{equation*}
$$

From Equs. (5), (6) and (7), we have

$$
\begin{equation*}
-\frac{\mathrm{d}\left[\mathrm{Br}_{2}\right]}{\mathrm{dt}}=\frac{2 \mathrm{k}_{1} \mathrm{k}_{2}[\mathrm{Qu}]^{2}}{\mathrm{k}_{-1}+\mathrm{k}_{2}\left[\mathrm{Q} \mathrm{Br}_{2}\right]} \tag{8}
\end{equation*}
$$

The reaction has been studied at the very low concentration of bromine (i.e. $15.52 \times 10^{-3} \mathrm{M}$ ) and hence if the $[\mathrm{Qu}]=\left[\mathrm{Br}_{2}\right]$ or the concentration of quinoline is slightly greater than the concentration of bromine, the values of $k_{2,2}[Q u]$ and $k_{-1}$ in the denominator of equ. 8. Can be neglected as the values are negligible and equ: (8) behaves as a third order expression-two with respect to quinoline and one with respect to bromine. When the concentration of quinoline is very lange as compared to bromine so that $\mathrm{k}_{2}[\mathrm{Qu}] \gg \mathrm{k}_{-1}$. The value of $\mathrm{k}_{-1}$ in denominator of equ. 8. Can be neglected, as it is negligible and the equ. (8). Can be written in the form.

$$
\begin{align*}
& -\frac{\mathrm{d}\left[\mathrm{Br}_{2}\right]}{\mathrm{dt}}=\frac{2 \mathrm{k}_{1} \mathrm{k}_{2}[\mathrm{Qu}]^{2}\left[\mathrm{Br}_{2}\right]}{\mathrm{k}_{2}[\mathrm{Qu}]}  \tag{9}\\
& -\frac{\mathrm{d}\left[\mathrm{Br}_{2}\right]}{\mathrm{dt}}=2 \mathrm{k}_{1}[\mathrm{Qu}]\left[\mathrm{Br}_{2}\right] \tag{10}
\end{align*}
$$

which is a second order expression, one with respect of each of the reactants.

## Acknowledgement

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# Effect of algae and phosphorus on the formation of aminoacids 

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#### Abstract

Blue green algae Anabaena and Tolypothrix fix atmospheric nitrogen, along with carbon saving. The efficiency of nitrogen fixation by algae is more when the systems are phosphated with TA FA or GERMAN basic slag. The amount and number of aminoacids is mo:e in the systems which are inoculated with algae, more in light than in the dark.


## Introduction

Since the discovery of phosphorus by Brandt ${ }^{1}$ in 1969 has lead the scientist to make intensive studies, because of the unique and paramount importarce of this element and its compounds to agriculture, biochemical processes and industry.

As in all living organisms, compounds containing phosphorus play important roles in nearly all phases of metabolism particularly in energy transformation reactions Krauss ${ }^{2}$ has emphasised the need of phosphorus for the optimum growth of algae. Free amino acids are rapidly degraded in soil their amount, however, may quickly increase after available carbohydrates and inorganic nitrogen have been supplied. According to several workers ${ }^{\text {s, }}$, the turnover of protein in soil depends on the nature of the organic matter being degraded and the environmental conditions. Bremner ${ }^{5}$ found however that organic matter tended to attain a constant protein and aminoacids composition. According to Watanabe ${ }^{6}$, Fowden ${ }^{7}$ and Williams and Burris ${ }^{8}$, aminoacids like aspartic, glutamic acid, alanine, glycine, valine serine and phenyl alanine are present in all the blue green algae. Dhar and Roy ${ }^{9}$ have also shown the formation of aminoacids by exposing normal solution of nitrate and glucose with titania. Simonart and Peeters ${ }^{10}$ studied the aminoacids of forest soils, leaf moulds and green house soil and reported that aspartic acid, glutamic acid, serine and alanine occured in all samples. Recently Payne et al ${ }^{11}$ noted that aqueous soil extraci concentrated by freeze drying yielded ninhydrin positive spots on paper.

In the present investigation the role of algae and phosphorus on the formation of amino acids in soil has been studied.

## Experimental

The experiments were carried out with the Institute soil. All the ingredients viz: soil, phosphates (Tata or German Basic Slag) and carbonaceous material (wheat straw) used in these experiments were sieved to 100 mesh .200 gms of soil weie taken in enamelled dishes. To this the required carbonaceous materials and phosphates were introduced as $1 \%$ carbon and $0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ respectively on an oven dry basis and the contents were thoroughly mixed in a big pestle and mortar and

[^1]the whole mixture was made as uniform as possible. Then wherever required the soil was inoculated with algae i.e. Chlorella vulgaris, Anabaena naviculoides and Tolypothrix tenuis. Two similar sets were arranged side by side, one of which was exposed to light under an electric bulb and the other was covered with a thick black cloth. In all the dishes the moisture content was maintained at $40 \%$ level throughout the experiment and the contents were stirred daily. At regular intervals of time, composite s.imples were taken and were analysed for total carbon, total nitrogen, available $\mathrm{P}_{21} \mathrm{O}_{5}$ and aminoacids. The quantitative estimation of aminoacids was made by Harding and Mclean ${ }^{12}$ procedure and qualitative detection by two dimensional paper chromatography. Bacteria free unialgal cultures of the above three algac were obtained from Dr. A. K. Mitra, an eminent algologist of the University of Allahabad.

The following abbreviations for the aminoacids studied have been used. Al for alanine, Va for valine, Asp for aspartic acid, Leu for leucine, Gly for Glycine, Glu for Glutamic acid, Ly for Lysine, Ar for Arginine and Hi for Histidine. The experimental observation are recorded in the following pages.

## Percentage analysis of the Ingredients

|  | Soil | Tata Basic <br> slag | German Basic <br> slag | Wheat straw |
| :--- | :---: | :---: | :---: | :---: |
| HCl insoluble | 78.326 | 15.7485 | 11.3978 | 5.126 |
| Sesquioxide | 9.2478 | - | - | 1.4154 |
| $\mathrm{Fe} \mathrm{O}_{3}$ | $4 \cdot 167$ | 15.4 .976 | 14.6893 | 0.6064 |
| CaO | 1.8765 | 38.6778 | 42.2897 | 0.8442 |
| $\mathrm{~K}_{2} \mathrm{O}$ | 1.1537 | 0.6337 | traces | 0.8036 |
| MgO | 1.5198 | 4.8247 | 4.9434 | 0.4078 |
| Total $\mathrm{P}_{2} \mathrm{O}_{5}$ | 0.2108 | 7.468 | 17.2614 | 0.6038 |
| ${\mathrm{Av} . \mathrm{P}_{2} \mathrm{O}_{5}}^{\text {Total carbon }}$ | 0.0136 | 4.1024 | 9.8876 | - |
| Total nitrogen | 0.7126 | - | - | 38.276 |

Table 1
200 gms Soil

| Period of <br> exposure <br> in days | Total <br> carbon <br> gm. | Total <br> nitrogen <br> gm. | Efficiency | Aminoacids <br> identified <br> chromatogra- <br> phically | Amount of <br> aminoacids <br> with respect <br> to glycine |
| :---: | :---: | :---: | :---: | :---: | :---: |


| LIGHT |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 0 | 1.4252 | 0.1396 |  | Al., Va., Glu. | ... |
| 60 | 1.3527 | $0 \cdot 1405$ | $12 \cdot 4$ * | Al., Va., Glu. |  |
| 120 | 1.3218 | 0.1408 | 11.6 | Al., Va., Glu., Ly. | 0.0309 |
| 180 | $1 \cdot 2757$ | 0.1412 | $10 \cdot 7$ | Al., Va., Gly., Ly., Ar. | 0.0348 |
| 240 | $1 \cdot 2042$ | $0 \cdot 1417$ | 9.5 | Al., Va., Glu., Ly., Ar. | 0.0334 |
| DARK |  |  |  |  |  |
| 0 | $1 \cdot 4252$ | $0 \cdot 13960$ |  | Al., Va., Glu. | $\cdots$ |
| 60 | 1.3765 | $0 \cdot 14000$ | $8 \cdot 2$ | Al., Va., Glu. |  |
| 120 | 1.3490 | $0 \cdot 14021$ | $8 \cdot 0$ | Al., Va., Glu. |  |
| 180 | 1.3186 | $0 \cdot 14040$ | $7 \cdot 5$ | Al., Va., Glu. Ly. | $0 \cdot 0271$ |
| 240 | $1.2 ヶ 64$ | $0 \cdot 14060$ |  | Al., Va., Glu., Ly. | 0.0276 |

Table 2
200 gms Soil + Chlorella vulgaris

| Period of <br> exposure <br> in days | Total <br> carbon <br> gm. | Total <br> nitrogen <br> gm. |  | Aminoacids <br> identified <br> chromatogra- <br> phically | Amount of <br> aminoacids <br> with respect <br> to glycine |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  | LIGHT |  |  |

Table 3
200 gms Soil + Anabaena naviculoides
LIGHT


Table 4
200 gms Soil + Tolypothrix tenuis

## LIGHT

| 0 | 1.4252 | $0 \cdot 1396$ |  | Al., Va., Glu. |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 60 | 1.3736 | 0.1412 | 31.0 | Al., Va., Glu., Ly |  |
| 120 | $1 \cdot 3504$ | 0.1417 | 28.0 | Al., Va., Gly., Ly. | 0.0335 |
| 180 | $1 \cdot 3126$ | $0 \cdot 1423$ | 2.0 | Al., Va., Glu., Ly. | $0 \cdot 0382$ |
| 240 | 1.2776 | $0 \cdot 1425$ | $19 \cdot 6$ | Al., Va., Glu., Ly | 0.0364 |
|  | DARK |  |  |  |  |
| 0 | $1 \cdot 4252$ | 0.13960 |  | Al., Va., Glu. | . |
| 60 | $1 \cdot 3755$ | $0 \cdot 14003$ | 8.4 | Al., Va., Glu. |  |
| 120 | $1 \cdot 3490$ | $0 \cdot 14019$ | 7.7 | Al., Va., Glu., Ly. |  |
| 180 | 1.3180 | $0 \cdot 14042$ | $7 \cdot 6$ | Al., Va., Glu., Ly. | 0.0271 0.0275 |
| 240 | 1.2854 | $0 \cdot 14062$ | $7 \cdot 2$ | Al., Va., Glu., Asp | 0.0275 |

Table 5
200 gms Soil $+1 \% \mathrm{C}$ as Wheat Straw

| Period of exposure in days | Total carbon gm. | Total nitrogen gm. | Efficiency | Aminoacids identified chromatographically | Amount of aminoacids with respcct to glycine |
| :---: | :---: | :---: | :---: | :---: | :---: |
| LIGHT |  |  |  |  |  |
| 0 | $3 \cdot 4252$ | $0 \cdot 1729$ |  | Al., Va., Glu. |  |
| 60 | 3.0584 | $0 \cdot 1833$ | 28.3 | Al., Va., Glu., Leu. | 1.7946 |
| 120 | $2 \cdot 8530$ | $0 \cdot 1887$ | $27 \cdot 6$ | Al., Va., Glu., Leu., | , Asp. $\quad 2 \cdot 4632$ |
| 180 | $2 \cdot 6779$ | $0 \cdot 1935$ | 27.5 | Al., Va., Glu., Asp., | , Leu.,Ar. $2 \cdot 7289$ |
| 240 | $2 \cdot 4709$ | 0.1926 | $27 \cdot 3$ | Al., Va., Glu., Asp. | $2 \cdot 6437$ |
| DARK |  |  |  |  |  |
| 0 | $3 \cdot 4252$ | 0.1729 |  | Al., Va., Glu. |  |
| 60 | 3.1404 | $0 \cdot 1774$ | $15 \cdot 8$ | Al., Glu., Va. | $1 \cdot 0721$ |
| 120 | $2 \cdot 9837$ | $0 \cdot 1797$ | $15 \cdot 4$ | Al., Glu., Va., Leu. | $1 \cdot 4897$ |
| 180 | $2 \cdot 8139$ | $0 \cdot 1821$ | $15 \cdot 3$ | Al., Glu., Va., Leu., | Asp. 1.9817 |
| 240 | $2 \cdot 6415$ | $0 \cdot 1818$ |  | Al., Va., Glu., Leu., | Asp., Ar. 1-8763 |
| Table 6 |  |  |  |  |  |
| 200 gms. Soil $+1 \% \mathrm{C}$ as Wheat Straw + Chlorella |  |  |  |  |  |
| LIGHT |  |  |  |  |  |
| 0 | 3.4252 | $0 \cdot 1729$ |  | Al., Glu., Va. |  |
| 60 | 3.1412 | $0 \cdot 1835$ | $37 \cdot 3$ | Al., Va., Glu., Leu., | Pro. 1.82 .5 |
| 120 | 2.9852 | 0.1891 | $36 \cdot 5$ | Al., Glu., Va., Leu., | Pro.,Asp. 2-4867 |
| 180 | 2.8074 | 0.1939 | $33 \cdot 9$ | Al., Va., Glu., Leu. Asp., Ar. | $\text { Pro., } \quad 2.7469$ |
| 240 | $2 \cdot 6199$ | $0 \cdot 1923$ |  | Al., Va., Glu., Leu., Asp., Ar. | Pro., 2.7231 |
| DARK |  |  |  |  |  |
| 0 | $3 \cdot 4252$ | $0 \cdot 1729$ |  | Al., Glu., Va. |  |
| 60 | $3 \cdot 1405$ | $0 \cdot 1774$ | $15 \cdot 8$ | Al., Glu., Va., Leu. | $1 \cdot 0722$ |
| 120 | $2 \cdot 9836$ | 0.1797 | $14 \cdot 4$ | Al., Va., Glu., Leu., | Al., Asp. $1 \cdot 4896$ |
| 180 | $2 \cdot 8140$ | $0 \cdot 1821$ | $15 \cdot 3$ | Al., Va., Glu., Leu., | , Asp., 1.9818 |
| 240 | $2 \cdot 6414$ | $0 \cdot 1819$ | ... | Al., Va., Glu., Leu., | Asp.,Ar. 1-8763 |

Table 7
200 gms. Soil $+1 \%$ C as Wheat Straw + Anabaena

## LIGHT

| 0 | 3.4252 | 0.1729 | $\ldots$. | Al., Va., Glu. | $\ldots$ |
| ---: | ---: | ---: | ---: | :--- | ---: |
| 60 | 3.1385 | 0.1860 | 45.6 | Al., Glu., Va., Se., Leu. | 2.0678 |
| 120 | 2.9841 | 0.1923 | 43.9 | Al., Glu., Va., Se,, Leu. | 2.7423 |
| 180 | 2.8063 | 0.1998 | 43.4 | Al. Va., Glu., Se., Leu., Asp., 2.9537 |  |
|  |  |  |  | Ar. |  |
| 240 | 2.6188 | 0.1986 | 42.2 | Al., Va., Glu., Se., Leu., Asp. 2.8234 |  |


| Period of <br> exposure <br> in days | Total <br> carbon <br> gm. | Total <br> nitrogen <br> gn. | Efficiency | Aminoacids <br> identified <br> chromatogra- <br> phically | Amount of <br> aminoacids <br> with respect <br> to glycine |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | DARK |  |  |  |  |
| 0 | 3.4252 | 0.1729 | $\ldots$ | Al., Glu., Va. |  |  |

Table 8
200 gms. Soil $+1 \% \mathrm{C}$ as Wheat Straw + Tolypothrix
LIGHT

| 0 | $3 \cdot 4252$ | $0 \cdot 1729$ |  | Al., Va., Glu. |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 60 | $3 \cdot 1398$ | 0.1867 | $48 \cdot 3$ | Al., Glu., Va., Th. | $2 \cdot 3467$ |
| 120 | $2 \cdot 984.5$ | $0 \cdot 1935$ | 46.7 | Al., Va., Glu., Asp., | $2 \cdot 9813$ |
| 180 | $2 \cdot 8069$ | $0 \cdot 2017$ | $46 \cdot 5$ | Al., Va., Glu., Asp., | 3.1294 |
| 240 | $2 \cdot 6196$ | $0 \cdot 2008$ |  | Al., Va., Glu., Th., Ar. | $2 \cdot 9672$ |
| DARK |  |  |  |  |  |
| 0 | $3 \cdot 4252$ | 0.1729 |  | Al., Va., Glu. |  |
| 60 | 3.1404 | 0.1774 | $15 \cdot 8$ | Al., Glu., Va., Leu. | 1.07898 |
| 120 | $2 \cdot 9836$ | 0.1797 | $15 \cdot 4$ | Al., Glu., Va., Leu. | 1.48818 |
| 180 | 2.8140 | $0 \cdot 1821$ | $15 \cdot 3$ | Al., Glu., Va., Leu. | 1.9818 |
| 240 | $2 \cdot 6416$ | $0 \cdot 1820$ |  | Al,, Va., Glu. | $1 \cdot 8764$ |

Table 9
200 gms. Soil $+1 \% \mathrm{C}$ as Wheat Straw $+0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ as Tata Basic Slag

## LIGHT



Table 10
200 gms. Soil $+1 \% \mathrm{C}$ as Wheat Straw $+0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ as T. B. S. + Chlorella

| Period of exposure in days | Total carbon gm. | Total nitrogen gm. | Efficiency | Aminoacids identified chromatographically | Amount of aminoacids with respect to glycine |
| :---: | :---: | :---: | :---: | :---: | :---: |
| LIGHT |  |  |  |  |  |
| 0 | $3 \cdot 4252$ | 0.1729 |  | Al., Va., Glu. |  |
| 60 | $2 \cdot 9684$ | $0 \cdot 1970$ | $52 \cdot 7$ | Al., Va., Glu., Pro. | 2.5984 |
| 120 | $2 \cdot 7160$ | 0.2092 | $51 \cdot 1$ | Al., Va., Glu., Pro., Leu. | 3.7818 |
| 180 | 2.5374 | $0 \cdot 2182$ | $51 \cdot 0$ | Al., Va., Glu., Pro., Leu., | , Hi., 45103 |
|  |  |  |  | Ar., Asp., Ar. |  |
| 240 | $2 \cdot 4173$ | 0.2169 | ... | Al., Va., Glu., Pro., Leu., Ar., Asp. | Hi., 4.2443 |

## DARK

| 0 | $3 \cdot 4252$ | 0.1729 | $\ldots$ | Al., Va., Glu. |  |
| ---: | :---: | :---: | :---: | :--- | ---: |
| 60 | 3.3904 | 0.1830 | $26 \cdot 1$ | Al., Va., Glu. | 1.6689 |
| 120 | $2 \cdot 84343$ | 0.1877 | $25 \cdot 4$ | Al., Glu., Va., Leu. | $2 \cdot 4369$ |
| 180 | $2 \cdot 63732$ | 0.1911 | $23 \cdot 0$ | Al., Va., Glu., Leu., Asp.,Hi. $2 \cdot 8187$ |  |
| 240 | $2 \cdot 48547$ | 0.1902 | $\ldots$ | Al., Va., Glu., Leu., Asp., Ar. $2 \cdot 7495$ |  |

Table 11
200 gms. Soil $+1 \% \mathrm{C}$ as Wheat Straw $+0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ as T. B. S. + Anabaena

## LIGHT

| 0 | $3 \cdot 4252$ | 0.1729 |  | Al., Va., Glu. |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 60 | $2 \cdot 9604$ | $0 \cdot 1985$ | $55 \cdot 0$ | Al. Va., Glu., Se. | $2 \cdot 8612$ |
| 120 | $2 \cdot 7152$ | 0.2110 | $53 \cdot 6$ | Al., Glu., Va., Se., Leu. | 4.0391 |
| 180 | $2 \cdot 5356$ | 0.2191 | $51 \cdot 9$ | Al., Va., Glu., Se., Leu., Hi., | $4 \cdot 8326$ |
| 240 | $2 \cdot 3866$ | 0.2183 |  | Asp., Ar. <br> Al., Va., Glu., Se., Leu., Hi., Asp., Ar. | 4.5349 |
| DARK |  |  |  |  |  |
| 0 | $3 \cdot 4252$ | 0.1729 |  | Al., Va., Glu. |  |
| 60 | $3 \cdot 1392$ | $0 \cdot 1830$ | $26 \cdot 1$ | Al., Va., Glu. | 1.6696 |
| 120 | $2 \cdot 84.43$ | $0 \cdot 1877$ | $25 \cdot 4$ | Al., Va.. Glu., Leu. | $2 \cdot 4358$ |
| 180 | $2 \cdot 63738$ | $0 \cdot 1911$ | $23 \cdot 0$ | Al., Va., Glu., Leu., Asp.,Hi. | $2 \cdot 8189$ |
| 240 | $2 \cdot 48541$ | $0 \cdot 1907$ | ... | Al., Va., Glu., Leu., Asp.. Hi., Ar. | $2 \cdot 7498$ |

Table 12
200 gms . Soil $+1 \% \mathrm{C}$ as Wheat Straw $+0 \cdot 5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ as T. B. S. + Tolypothrix

| Period of <br> exposure <br> in days | Total <br> carbon <br> gm. | Total <br> nitrogen <br> gm. | Efficiency | Aminoacids <br> identified <br> chromatogra- <br> phically | Amount of <br> aminoacids <br> with respect <br> to glycine |
| :---: | :---: | :---: | :---: | :---: | :---: |

## LIGHT

| 0 | $3 \cdot 4252$ | $0 \cdot 1729$ |  | Al, Glu., Va. |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 60 | $2 \cdot 9629$ | $0 \cdot 1992$ | $56 \cdot 8$ | Al., Glu., Va., Th. | $3 \cdot 1327$ |
| 120 | $2 \cdot 7159$ | $0 \cdot 2112$ | $54 \cdot 4$ | Al., Glu., Va., Th., Leu. | $4 \cdot 3672$ |
| 180 | $2 \cdot 5366$ | $0 \cdot 2200$ | 53.0 | Al., Glu., Va., Th., Leu., | 5•1853 |
| 240 | $2 \cdot 3580$ | 0.2194 | $\ldots$ | Asp., Ar., Hi. <br> Al., Glu., Va., Th., Leu., Asp., Ar., Hi. | 4.9401 |

DARK

| 0 | $3 \cdot 4252$ | $0 \cdot 1729$ |  | Al., Glu., Va. |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 60 | $3 \cdot 0392$ | $0 \cdot 1830$ | $26 \cdot 1$ | Al., Va., Glu., Leu. | 1.6698 |
| 120 | $2 \cdot 4434$ | $0 \cdot 1877$ | $25 \cdot 5$ | Al., Glu., Va., Leu. | 2.4369 |
| 180 | $2 \cdot 63736$ | 0.1911 | $23 \cdot 0$ | Al., Glu., Va., Leu., Hi. | 2.8186 |
| 240 | $2 \cdot 4854$ | 0.1907 |  | Al., Glu., Va., Leu., Hi, Ar. | $2 \cdot 7493$ |

Table 13
200 gms. Soil $+1 \% \mathrm{C}$ as Wheat Straw $+0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ as German Basic Slag LIGHT

| 0 | $3 \cdot 4252$ | 0.1729 |  | Al., Va., Glu. |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 60 | $2 \cdot 8340$ | $0 \cdot 1982$ | $42 \cdot 7$ | Al., Glu., Va., Leu. | $2 \cdot 7467$ |
| 120 | $2 \cdot 6559$ | 0.2051 | $41 \cdot 8$ | Al., Glu., Va., Leu., Asp. | 3.8923 |
| 180 | $2 \cdot 3974$ | 0.2157 | $41 \cdot 6$ | Al., Glu., Va., Leu., Asp., Ar., Hi. | $4 \cdot 5871$ |
| 240 | 2.2784 | 0.2146 |  | Al., Va., Glu., Leu., Asp., Ar., Hi. | $4 \cdot 4108$ |
| DARK |  |  |  |  |  |
| 0 | $3 \cdot 4252$ | $0 \cdot 1729$ |  | Al., Glu., Va. |  |
| 60 | $3 \cdot 0376$ | $0 \cdot 1831$ | $26 \cdot 2$ | Al., Va., Glu., Leu. | 1.7224 2.4178 |
| - 120 | $2 \cdot 8426$ | $0 \cdot 1877$ | $25 \cdot 4$ | Al., Va., Glu., Leu., Hi. | 2.4178 |
| 180 | $2 \cdot 6361$ | $0 \cdot 1913$ | $23 \cdot 3$ | Al. Va., Glu., Leu., Hi. | $2 \cdot 9119$ |
| 240 | $2 \cdot 4980$ | $0 \cdot 1904$ | ... | Al., Va., Glu., Leu., Hi., Asp | $2 \cdot 8274$ |

Table 14
200 gms. Soil $+1 \% \mathrm{C}$ as Wheat Straw $+0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ as $\mathrm{G} . \mathrm{B} . \mathrm{S} .+$ Chlorella
LIGHT

| 0 | 3•4252 | $0 \cdot 1729$ |  | Al., Va., Glu. |
| :---: | :---: | :---: | :---: | :---: |
| 60 | $2 \cdot 9600$ | 0. 1985 | $55 \cdot 0$ | Al. Va., Glu., Pro. |
| 120 | $2 \cdot 7189$ | $0 \cdot 2093$ | $51 \cdot 5$ | Al., Va., Glu., Pro., Leu.,Asp. ${ }^{\text {3.93153 }}$ |
| 180 | 2.5298 | 0.2187 | $51 \cdot 1$ | Al., Va., Glu., Pro., Leu., <br> Asp., Hi. |
| 240 | $2 \cdot 3840$ | 0.2179 | $\ldots$ | Al., Glu., Va., Pro., Leu., 4.4549 Asp., Hi., Ar. |


| Period of <br> exposure <br> in days | Total <br> carbon <br> gm. | Total <br> nitrogen <br> gm. | Efficiency | Aminnacids <br> identified <br> chromatogra- <br> phically | Amount of <br> aminoacids <br> with respect <br> to glycine |
| :---: | :---: | :---: | :---: | :---: | :---: |

## DARK



Table 15
200 gms. Soil $+1 \% \mathrm{C}$ as Wheat Straw $+0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ as G. B. S. + Anabaena LIGHT

| 0 | $3 \cdot 4252$ | $0 \cdot 1729$ |  | Al., Glu., Va. |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 60 | $2 \cdot 9581$ | $0 \cdot 1996$ | 57-1 | Al., Glu., Va., Se. | 3.2659 |
| 120 | $2 \cdot 7163$ | $0 \cdot 2117$ | $54 \cdot 6$ | Al., Glu., Va., Se., Leu.,Asp. | $4 \cdot 4165$ |
| 180 | $2 \cdot 5280$ | 0.2204 | $52 \cdot 9$ | Al., Va., Glu., Se., Lea., Asp., Hi. | 4.8724 |
| 240 | $2 \cdot 3828$ | 0.2197 |  | Al., Va., Glu., Se., Leu., Asp., Hi. | $4 \cdot 689$ |
| DARK |  |  |  |  |  |
| 0 | $3 \cdot 4252$ | 0.1729 |  | Al., Glu., Va. |  |
| 60 | 3.03762 | $0 \cdot 1831$ | $26 \cdot 2$ | Al., Glu., Va., Leu. | 1.7226 |
| 120 | $2 \cdot 84262$ | 0.1877 | $25 \cdot 4$ | Al., Glu., Va., Leu. | $2 \cdot 4181$ |
| 180 | $2 \cdot 6304$ | $0 \cdot 1913$ | $23 \cdot 3$ | Al., Glu., Va., Leu., Asp. | $2 \cdot 9119$ |
| 240 | $2 \cdot 4986$ | $0 \cdot 1904$ | ... | Al., Glu., Va., Leu., Asp., Hi. | $2 \cdot 8273$ |

Table 16
200 gms. Soil $+1 \% \mathrm{C}$ as Wheat $\mathrm{Stam}^{2}+0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ as G. B. S. + Tolypothrix
LIGHT

| 0 | $3 \cdot 4252$ | 0.1729 |  | Al., Glu., Va. |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 60 | $2 \cdot 9584$ | 0.2004 | 58.8 | Al., Glu., Va., Leu., Th. | $3 \cdot 3729$ |
| 120 | $2 \cdot 7180$ | 0.2135 | $57 \cdot 4$ | Al., Glu., Va., Leu., 'Th., Asp., Ar. | 4.5237 |
| 180 | 2.5289 | 0.2210 | $53 \cdot 6$ | Al., Va., Glu.; Leu., Th., Asp., Ar., Hi. | $4 \cdot 9635$ |
| 240 | $2 \cdot 3837$ | 0.2205 | $\cdots$ | Al., Glu., Va., Leu., Th., Asp., Ar., Hi. | $4 \bullet 6812$ |
| DARK |  |  |  |  |  |
| 0 | 3.4252 | $0 \cdot 1729$ |  | Al., Glu., Va. |  |
| 60 | $3 \cdot 03761$ | $0 \cdot 1831$ | $26 \cdot 2$ | Al., Glu., Va., Leu. | 1.7225 |
| 120 | 2.84260 | 0.1877 | $25 \cdot 4$ | Al., Glu., Va., Leu., Asp. | $2 \cdot 4179$ |
| 180 | $2 \cdot 6.612$ | $0 \cdot 1913$ | $23 \cdot 3$ | Al., Glu., Va., Leu., Asp. | $2 \cdot 9121$ |
| 240 | $2 \cdot 49801$ | 0.1904 | ... | Al., Glu., Va., Leu., Asp., | Hi. $2 \cdot 8272$ |

## Discussion

It is observed that there is an appreciable oxidation of carbon from wheat straw, when mixed with soil and allowed to undergo slow oxidation in air at the laboratory temperature, more in the light than in the dark sets. There is concomitant increase in the number and amount of aminoacids formed along with nitrogen fixation, which is also greater in the sets exposed to light than in the dark. These observations show that light energy is helpful in nitrogen fixation and subsequent utilization of the fixed nitrogen in the synthesis of amınoacids.

It is very interesting to note that the efficiency of nitrogen fixation i.e. the amount of nitrogen fixed in milligram per gram of carbon oxidised is always greater in the sets exposed to light than those kept in the dark and the efficiency gradully decreases, the reasons are clear from the following consideration. In our previous paper ${ }^{13}$ it is shown, that in soils, the process of nitrogen fixation is always opposed by ammonification and nitrification. Thus on one hand, nitrogen fixed on the surface of the soil undergoes various changes aided by light absorption to form aminoacids and proteins and on the other hand these aminoacids and proteins formed in the soil or sand undergo ammonification and nitrification which are accelerated by light absorption and form nitrite and nitrate in the following manner :-

$$
\text { Proteins } \longrightarrow \text { Aminoacids } \xrightarrow{+\mathrm{O}_{2}} \mathrm{NH}_{4} \text { compounds } \xrightarrow{+\mathrm{O}_{2}} \mathrm{NO}_{2} \xrightarrow[\longrightarrow]{+\mathrm{O}_{2}} \mathrm{NO}_{3}
$$

In these series of reactions the unstable ammonium nitrite is formed which undergoes decomposition liberating energy and nitrogen gas as in the following equation :-

$$
\mathrm{NH}_{4} \mathrm{NO}_{2} \rightarrow \mathrm{~N}_{2}+2 \mathrm{H}_{2} \mathrm{O}+718 \mathrm{~K} . \text { cal. }
$$

But the addition of phosphate i.e. Tata or German basic slag to the systems, form stable phosphoproteins or phospho ligno proteins which resist the ammonification and nitrification processes. It is clear from the results that the amount and number of aminoacids is greater in phosphated than in the unphosphated systems. It is interesting to record that in the sets inoculated with algae some saving in carbon and a small increase in nitrogen was always found in light though in the dark, there was no significant difference in the carbon and nitrogen status of the systems. The order of carbon saving was Chlorella $>$ Tolypothrix $>$ Anabaena. In sets inoculated with chlorella, there is a slight increase in nitrogen compared with those which have not been inoculated. It is interesting to note that though chlorella is not a fixer of nitrogen, even then it shows a small fixation of nitrogen in light. The probable reason of this seems to be that in soils chlorella lives in symbiosis with nitrogen fixing organism azotobacter ${ }^{14}$ and during this process it supplies the azotobacter with carbohydrates and the bacteria in turn fixes nitrogen.

In the sets inoculated with tolypothrix and Anabaena, there is appreciable increse in the nitrogen content over that of control sets as well as in sets inoculated with chlorella, thereby having better efficiencey, the increase being greater in the case of tolypothrix than in that of anabaena. This is in agreement with the general observation made by various algologists that tolypothrix is a better fixer of nitrogen than anabaena.

The systems inoculated with algae contain greater amount and number of aminoacids than is contained in the uninoculated systems. This is probably due to decomposition and oxidation of algal bodies which increases the amount and
number of aminoacids. It has been generally observed that chlorella favours the formation of proline, whil: in systems having anabaena and tolypothrix serine and threonine have generally been detected. The amount of aminoacids is always larger in sets inoculated with anabaena and tolypothrix than in the sets containing chlorella. The posible reason for this seems to be that anabaena and tclypoth:ix fix atmospheric nitrogen along with carbon saving, the fixed nitrogen is utilized in the synthesis of aminoacids and some of the aminoacids are excreted in the medium. Fowden ${ }^{7}$, Watanabe ${ }^{6}$ and Williams and Burris, ${ }^{8}$ found many free aminoacid present in algae. Rothwell and Frederick ${ }^{4}$ observed aminoacids during the decomposition of alialfa and corn stover in soil at different temperatures.

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# Oxidation of $\alpha$-Hydroxy acids by Vanadium (V) 

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#### Abstract

The stoichiometry of the oxidation of $\alpha$-hydroxy acids by vanadium (V) under different conditions of mineral acid concentrations has been establisbed potentiometrically using platinum and meicury coated platinum electrodes. At low $\mathrm{H}_{2} \mathrm{SO}_{4}$ concentrations 2 equivalents and at high $\mathrm{H}_{2} \mathrm{SO}_{4}$ concentrations 4 equivalents of the oxidant are used up per mole of the hydroxy acid, yielding corresponding aldehydes and carboxylic acids respectively. At intermediate concentrations of $\mathrm{H}_{2} \mathrm{SO}_{4}$ the lack of stoichiometry is due to two simulaneous competitive reactions leading to mixed products. These results ate explained on the basis of the standard oxidation potential data. Products in some cases have been isolated and identified.


The reactions have analytical significance.

In the oxidation of lactic acid by pericdic acid Brown ${ }^{1}$ realised all the successive oxidation stages (pyruvic acid, acetic acid, formic acid and finally $\mathrm{CO}_{2}$ ). In the kinetic studies of the oxidations of $\alpha$-hydroxy acids by different oxidising agents, other workers ${ }^{2-8}$ have assumed the oxidation products to be $\alpha$-keto acids, aldehydes or simple carboxylic acids depending upon the nature of the oxidant and the concentration of the mineral acid used.

In the present paper, the stoichiometry of the oxidation of lactic and mandelic acids with vanadium ( $V$ ) under varying concentrations of mineral acids has been determined potentiometrically in addition to exploring analytical prospects of such oxidations.

## Experimental

$\mathrm{M} / 100$ solutions of $\mathrm{NaVO}_{3}$ (Analar) were prepared, in different concentrations of $\mathrm{H}_{2} \mathrm{SO}_{4}$ and standardised with Mohr's salt solution using N-phenyl anthranilic acid as indicator. $\mathrm{M} / 20$ solutions of lactic and mandelic acids were prepared in freshly distilled air free water and standardised with carbonate free NaOH solutions.

For potentiometric investigations, 20 ml . of the acid vanadate $+x \mathrm{ml}$. of a-hydroxy acid, $x$ varying from 0 to 2, were taken in a series of bottles. After vigorous shaking, the raction mixtures were allowed to remain at $28^{\circ} \mathrm{C}$ for 24 hours to attain equilibrium. The indicator electrode, in the case of mandelic acid, was a bright platinum foil and in the case of lactic acid was a mercury coated thick
platinum wire. While in the case of lactic acid the bright platinum foil failed to detect any point of inflexion, in the case of mandelic acid both the electrodes worked satisfactorily, the mercury electrode proving better (vide fig. 3). E. m. f's were determined by immersing the electrode in each bottle, one by one, and combining it withi: saturated calomel electiode.

The ${ }^{i} \mathrm{Hg}$ coated Pt electrode was prepared by the method described by Strouts, Gilfillen and Wilson. ${ }^{7}$

The Pt electrode, prior to its use, was cleaned in chromic acid followed by treatment with ( N ) NaOH and subsequent washing with distilled water and glowin alcohol flame.


con

Table 1
20 ml . of $\mathrm{x}(\mathrm{M} / 100) \mathrm{NaVO}_{3}$ in $y(\mathrm{~N}) \mathrm{H}_{2} \mathrm{SO}_{4} \equiv \mathrm{z} \mathrm{ml}$. of $\mathrm{M} / 20 \alpha$-hydroxy acid. Temperature- $28^{\circ} \mathrm{C}$

| a-hydroxy acid | x | y | z | moles of $\mathrm{NaVO}_{3}$ per mole of <br> the $\alpha$-hydroxy acid |
| :---: | :---: | :---: | :---: | :---: |
| Lactic acid $0.95(\mathrm{M} / 20)$ | 0.901 | 0.5 | 1.9 | 2.00 |
|  | 0.901 | 1 | 1.9 | 2.00 |
|  | 0.901 | 4 | 1.9 | 2.00 |
|  | 0.950 | 8 | 1.7 | 2.35 |
|  | 0.950 | 10 | 1.2 | 3.33 |
|  | 0.950 | 14 | 1.1 | 3.63 |
| Mandelic acid $(\mathrm{M} / 20)$ | 0.960 | 18 | 1.0 | 4.00 |
| $\ldots$ | 0.98 | 0.1 | 1.8 | 2.1 |
|  | 0.98 | 0.5 | 1.8 | 2.1 |
|  | 0.98 | 1 | 1.8 | 2.1 |
|  | 0.98 | 4 | 1.4 | 2.7 |
|  | 0.98 | 8 | 1.1 | 3.45 |
|  | 0.98 | 18 | 1.0 | 3.92 |

Volumetric determinatzon of the o-hydroxy actds by oxidations with acid vanadate :
10 ml . of a solution containing a known weight of the a-hydroxy acid and 30 ml . of 0.95 ( $\mathrm{M} / 10$ ) $\mathrm{NaVO}_{3}$ were mixed to which sufficient amount of concentrated $\mathrm{H}_{2} \mathrm{SO}_{4}$ was acded so as to bring the acid concentration above 14 N . After thoroughly shaking and leaving the reaction mixture for 24 hours, the solutions were diluted to 100 ml . The excess of vanadate in solutions was back titrated with a standard Mohis salt soluti in using $N$ phenyl anthranilic acid as indicator. At the end point the colour change was from violet to green.

Table 2

| $\mathrm{H}_{2} \mathrm{SO}_{4}$ concn. | $\alpha$-hydroxy acid taken | Found | Error |
| :---: | :--- | :---: | :---: |
| 16 N | 45 mg. of lactic acid | 44.1 mg. | $2 \%$ |
| 18 N | 45 mg . of lactic acid | 44.4 mg. | $1.33 \%$ |
| 16 N | 76 mg. of mandelic acid | 74.95 mg. | $1 \cdot 38 \%$ |
| 18 N | 76 mg. of mandelic acid | 75.1 mg. | $1 \cdot 18 \%$ |

## Results and Discussions

Preliminary investigations as well as results in Table 1 indicate that both the rate and stoichiometries of the present oxidutions are mineral acid dependent, the oxidation of the mandelic acid being faster than that of the lactic acid. In the complete absence of mineral acids, no reaction occurs; at concentrations below $N / 2$, no inflexion or vague ones were located in the e. m. f. curves. At
$\mathrm{H}_{2} \mathrm{SO}_{4}$ concentrations $\mathrm{N} / 2$ to $4 \mathrm{~N}, 2.0$ moles, and above $14 \mathrm{~N}, 4 \cdot 0$ moles of vanadium (V) are consumed per mole of the hydroxy acid. At intermediate concentrations of $\mathrm{H}_{2} \mathrm{SO}_{4}$ the moles of vanadium ( V ) consumed per mole of the hydroxy acid are not integral, though the value increases and approaches a limiting one of $4 \cdot 0$, as the $\mathrm{H}_{2} \mathrm{SO}_{4}$ concentration increases. Assuming one equivalent reduction of vanadium (V) in each case, these results correspond to the following overall redox processes,

$$
\begin{array}{r}
\text { At bigh } \mathrm{H}_{2} \mathrm{SO}_{4} \text { conicn., } \mathrm{RCHOHCOOH}+4 \mathrm{VO}_{2}{ }^{+}+4 \mathrm{H}^{+}= \\
+4 \mathrm{COOH}+\mathrm{CO}_{2}+3 \mathrm{H}_{2} \mathrm{O} \\
+\quad+\quad .
\end{array}
$$

$$
\text { At low } \mathrm{H}_{2} \mathrm{SO}_{4} \text { concn., } \mathrm{RCHOHCOOH}+2 \mathrm{VO}_{2}^{+}+2 \mathrm{H}^{+}=\underset{\mathrm{RCOCOOH}+2 \mathrm{H}_{2} \mathrm{O}}{ }+2 \mathrm{VO}++\quad \text { + }
$$

$$
\text { or } \quad \mathrm{RCHOHCOOH}+2 \mathrm{VO}_{2}++2 \mathrm{H}^{+}=\mathrm{RCHO}+\mathrm{CO}_{2}+2 \mathrm{H}_{2} \mathrm{O}
$$

$$
\begin{equation*}
+2 \mathrm{VO}++ \tag{il}
\end{equation*}
$$

No significant change was noticed, except in the values of e. m. f., by sub,tiuting HCl or $\mathrm{HClO}_{4}$ for $\mathrm{H}_{2} \mathrm{SO}_{4}$. The use of $\mathrm{H}_{3} \mathrm{PO}_{4}$ was difficult because the vanadate solution tended to through a yelluw precipitate, presumably a hetercpoly acid.

At low acid-concentrations, the reaction mixture, in case of mandelic acid, soon smelt of benzaldehyde which was isolated in the end and characterised as its 2 : 4-dinitrophenyl hydrazone (m. p. $235^{\circ} \mathrm{C}$ ). In the case of lactic acid the acetaldehyde formed was isolated and characterised as its dimedone. derivative ( $\mathrm{m}, \mathrm{p} .140^{\circ} \mathrm{C}$ ). In the case of mandelic acid, at high $\mathrm{H}_{2} \mathrm{SO}_{4}$ concenuations the product, benzoic acid, was isolated, recrystallised and its m. p. checked. Any $n-\mathrm{keto}$ acid formed as an intermediate is, therefcre, ultimately, with lap e of time, converted to the aldehyde without affecting the stoichiometry of the redox process.

At intermediate acid concentrations ( 4 N to 14 N ) the lack of stoichiometry is presumed to be due to two simultaneous competitive reactions (i) and (ii) leading to mixed products. The fact that in no case more than one inflexion is no.iced in the e.m.f. curves leads to the conclusion that the same equilibrium constant covers both products.

Oxidation to formic acid (in case of lactic acid) and $\mathrm{CO}_{2}$, requiring 10 and 12 equivalents of the oxidant respectively, is not indicated even at acid concentrations greater than 14 N .

The acid dependence of the reaction rate is understood in terms of the existence of an equilibrium

$$
\mathrm{VO}_{2}{ }^{+}+\mathrm{H}_{3} \mathrm{O}^{+}=\mathrm{V}(\mathrm{OH})_{3}{ }^{++}
$$

followed by the attack of the doubly charged cation which is known to be a more powerful oxidiser than the $\mathrm{VO}_{2}{ }^{+}$itsell ${ }^{8}$. The range of redox process which can then occur is limited by thermodynamic considerations. On the basis of the potentials of the lactic acid-pyruvic acid couple ${ }^{\ominus} i \mathrm{E}^{\circ}=-0.20$ vol!) and acetalde-hyde-acetic acid couple ${ }^{10}$ ( $\mathrm{E}^{\circ}=0.118$ volt) it should be relatively easy to oxidise lactic acid to acetic acid than to pyruvic acid (or acetaldehyde). Yet both vanadate and dichromate ${ }^{5}$, which differ in their $\mathrm{E}^{\circ}$ values by 0.30 volt only, oxidise, at low acid concentrations, lactic acid to pyruvic acid (or acetaldehyde), oxidation to acetic acid occurring only at high $\mathrm{H}_{2} \mathrm{SO}_{4}$ concentrations. With respect to the
reduction of vanadate, larger decreases of free energy are expected in the reduction to $\mathrm{V}(\mathrm{IV})$ ( $\mathrm{E}^{\circ}=-1.000$ volt ) rather than to V (III) ( $\mathrm{E}^{\circ}=-0.68 \mathrm{volt}$ ). Moreover, comparing with other weak 2-equivalent oxidants like $\mathrm{I}_{2}$ ( $\mathrm{E}^{\circ}=-0 \cdot 53$ volt) or arsenic acid ( $\mathrm{E}^{\circ}=-0.559$ volt $)$, the $\mathrm{V}($ III $)-\mathrm{V}(\mathrm{IV})$ couple is not likely to be involved in organic oxidations.

The stoichiometric oxidations of the a-hydroxy acids at high $\mathrm{H}_{2} \mathrm{SO}_{4}$ concentrations raise the prospects of their analytical applications. An attempt has been made to evolve a method with an error of less than $2 \%$ (Table 2).

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# Effect of Direct and Indirect application of Organic Matter (phosphated and unphosphated) on the yield of wheat crop <br> By 

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The soil fertility is impoverised by continuous cropping. The capacity of the soil to supply nutrient elements to plants in required amounts and balanced proportions is commonly termed as "Soil fertility". It is governed by several factors. Of them, most important is the organic matter. As soil organic matter, in addition to amelioration of physical properties of soil, improves nutrient status of the soil, when it undergoes slow oxidation in the soil. A majority of soil scientists, have emphasised a need of replenishing the supply of organic matter in the soil. Karunakar, R. D. (1951) pointed out that organic matter adds nutrients to the soil. It has also been revealed by a number of soil scientists that organic matter chelates the nutrients available in the soil and thus checks the loss of available nutrients from soil.

In Indian soils organic matter is lower as compared to the soils of European countries. It ranges from 0.4 to $0.6 \%$. In view of the vital role of the organic matter in the land fertility and its general deficiency in our soils, one of the major problems of Indian Agriculture is to build up soil organic matter and to conserve it.

The sources of the organic matter are crop residues, green manures, animal refuse, weeds and tree leaves. Out of these, crop residues can not be utilised by transporting to another field, and green manuring is not always profitable due to usual corp rotations. The animal refuse, weeds, tree leaves have other uses like fuels than as manure. Composting is advocated for this purpose, yet composts have a doubtful utility in the building up the organic matter content of soil. It is because of heavy losses of nitrogen during composting process which is reported by many scientists mainly Albrecht and Poirst, Gilbert, Pember, Gerlach etc.

Dhar and co-workers have found that direct application of plant materials to the fields, instead of composting, is more beneficial to crops because the energy materials like carbohydrates, fats, lignin etc. when added to the soil, are slowly, oxidised. The energy liberated helps in fixing atmospheric nitrogen on the soil surface, forming $\mathrm{NH}_{3}$, Amino acids and other nitrogenous compounds.

A study of comparing organic matter used directly and after composting has been made by us in pots in the green house of the Institute under controlled conditions in the wheat.

## Experimental

Soil of top $9 "$ was taken, dried, powdered; cleaned from foreign materials mixed well before use. 6 kgs of this soil were taken in which organic marerials
e. g. Paddy straw, water hyacinth and weed (after thorough mixing, powdering, and sieving through 100 mesh ) were added @ 10 tons/acre. Phosphate was added as G. B. S. @ 50 lbs. $\mathrm{P}_{2} \mathrm{O}_{5}$ purcent. For indirect method composts (phosphated with G. B. S. containing $17.4 \% \mathrm{P}_{2} \mathrm{O}_{5}$ ) were already prepared. The rate of compost application was on nitrogen basis which was 100 lbs . per acie.

After thorough mixing the mixture was filled in pots and the system was kept moist with $20 \%$ moisture level by weight. In case of direct application the pots were exposed to sunlight for one month with mixing of materials at an interval ol a week, where as in composting the soil and composts were thoroughly mixed and exposed for 15 days with two mixings and $20 \%$ moisture level. Then the wheat crop was sown in both the systems on the same date.

## Scheme:

1. Soil alone
2. Soil + mixed leaves (weeds)
3. Soil + weed + G. B. S. @ $1 \% \mathrm{P}_{2} \mathrm{O}_{5}$.
4. Soil + weed + G.B. S. @ $0.5 \% \mathrm{P}_{2} \mathrm{O}_{8}$.
5. Soil + water hyacinth
6. Soil + water hyacinth + G. B. S. @ $1 \% \mathrm{P}_{2} \mathrm{O}_{5}$.
7. Soil $+, \quad, \quad+\quad, 0.5 \% \mathrm{P}_{2} \mathrm{O}_{8}$.
8. Soil + Paddy Straw
9. Soil $+,, \quad+\quad, \quad$ @ $1 \% \mathrm{P}_{2} \mathrm{O}_{5}$.
10. Soil $+, \quad, \quad+\quad, \quad 0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$.

Treatments were the same in direct and indirect applications.

| Replications | 4 |
| :--- | :--- |
| Total no. of pots | 40 |
| Design | Randomised Block Design |
| Date of sowing | 15 th Oct. 1968 |
| Date of harvesting | 28th April 1969 |
| Variety | K-68 |
| Plant population (per pot) | 5 |

Materials taken for the application :
Organic matters :

| Analysis | Water hyacinth | Paddy straw | Weed (Mixed) |
| :--- | :---: | :---: | :---: |
| Loss on ignition | 63.241 | 86.23 | 89.45 |
| Ash | 38.53 | 13.04 | 10.13 |
| $\mathrm{Fe}_{2} \mathrm{O}_{3}$ | 1.34 | 1.33 | 0.42 |
| $\mathrm{Al}_{2} \mathrm{O}_{3}$ | 0.089 | - | - |
| CaO | 2.263 | 0.66 | 0.74 |
| MgO | 0.785 | 0.169 | 0.708 |
| $\mathrm{~K}_{2} \mathrm{O}$ | 4.684 | 1.706 | 0.983 |
| $\mathrm{P}_{2} \mathrm{O}_{5}$ | 0.563 | 0.710 | 0.526 |
| Total carbon | 32.44 | 31.996 | 41.349 |
| Total nitrogen | 1.514 | 0.580 | 0.659 |
| C/N ratio | 21.56 | 55.17 | 63.61 |

Analysis of German Basic Slag which is used as phosphate

| Silica | 13.227 |
| :--- | ---: |
| $\mathrm{Fe}_{2} \mathrm{O}_{\mathrm{s}}$ | 16.136 |
| Total $\mathrm{P}_{2} \mathrm{O}_{5}$ | 17.853 |
| Available $\mathrm{P}_{2} \mathrm{O}_{5}$ | 8.997 |
| Total GaO | 41.132 |
| Total K O | 0.184 |
| Total MgO | 5.103 |

Analysis of compost prepared (phosphated and unphosphated)
Phosphated with G. B. S.
Analysis $\quad 1 \%$ level $0.5 \%$ level Unphosphated

Mixed weeds compost
Total carbon
Total nitrogen
$\mathrm{NH}_{5}-\mathrm{N}$
$\mathrm{NO}_{3}-\mathrm{N}$
Available $\mathrm{P}_{2} \mathrm{O}_{8}$

Total carbon
Total nitrogen
$\mathrm{NH}_{8}-\mathrm{N}$
$\mathrm{NO}_{3}-\mathrm{N}$
Available $\mathrm{P}_{2} \mathrm{O}_{8}$

| 12.93 | 16.04 | 25.01 |
| :--- | :---: | :---: |
| 1.68 | 1.49 | 1.16 |
| 0.062 | 0.043 | 0.0431 |
| 0.1214 | 0.102 | 0.0896 |
| 0.721 | 0.385 | 0.282 |

Paddy straw compost

| 9.75 | 11.03 | 19.68 |
| :--- | :---: | :---: |
| 1.58 | 1.36 | 0.93 |
| 0.084 | 0.046 | 0.271 |
| 0.139 | 0.963 | 0.072 |
| 0.649 | 0.341 | 0.2639 |

Water hyacinth compost

| Total carbon | 16.491 | 17.680 |  |
| :--- | ---: | :---: | :---: |
| Total nitrogen | 2.594 | 20.143 |  |
| $\mathrm{NH}_{3}-\mathrm{N}$ | 0.136 | 0.406 | 1.937 |
| $\mathrm{NO}_{3}-\mathrm{N}$ | 0.198 | 0.2013 | 0.0814 |
| Available $\mathrm{P}_{2} \mathrm{O}_{5}$ | 0.697 | 0.335 | 0.1030 |
|  |  |  |  |

Average Yield of wheat crop (in grams)

| Treatments | Direct (grain) | Indirect (grain) |
| :--- | :---: | :---: |
| Soil alone | $53 \cdot 1$ | $50 \cdot 6$ |
| Soil + mixed weed | $136 \cdot 3$ | $68 \cdot 5$ |
| Soil + mixed weed $+1 \% \mathrm{P}_{2} \mathrm{O}_{5}$ | $151 \cdot 3$ | $79 \cdot 3$ |
| Soil + mixed weed $+0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ | $146 \cdot 2$ | $72 \cdot 2$ |
| Soil + water hyacinth | $149 \cdot 6$ | $73 \cdot 8$ |
| Soil + water hyacinth $+1 \% \mathrm{P}_{2} \mathrm{O}_{5}$ | $167 \cdot 4$ | $89 \cdot 5$ |
| Soil + water hyacinth $+0.5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ | $152 \cdot 9$ | $76 \cdot 8$ |
| Soil + Paddy straw | $131 \cdot 4$ | $62 \cdot 0$ |
| Soil + Paddy straw $+1 \% \mathrm{P}_{2} \mathrm{O}_{5}$ | $146 \cdot 2$ | $70 \cdot 1$ |
| Soil + Paddy straw $+0 \cdot 5 \% \mathrm{P}_{2} \mathrm{O}_{5}$ | $139 \cdot 5$ | $67 \cdot 5$ |

## Discussion

An examination of the yield data manifests a marked increase in yield of wheat crop in all treatments in which organic materials such as paddy straw,
mixed weeds, water hyacinth have been added either directly or indirectly over control. Moreover, when organic matter is applied in conjunction with calcium phosphates German Basic Slag higher yields were obtained over the control as well as even the treatments in which only organic matter was added.

Dhar and his coworkers have established this that slow oxidation of organic matter fixes atmospheric nitrogen. Further they have also proved that addition of phosphates in conjunction with organic matter stabilises the fixed atmospheric nitrogen due to phosphorylation. From this fact, it can be stated that increase of wheat yield, in treatments in which organic matter and organic matter with phosphates were added, is due to fixation of atmospheric nitrogen as well as the availability of phosphate, calcium, Mg., potassium and to some extent trace elements due to dissolution of basic slag in the soil on account of liberating of organic acids chiefly carbonic acid owing to slow oxidation of organic matter in soil.

Direct application of organic matter (phosphated and unphosphated) is better than that of indirect (compost) because in case of direct application more nitrogen was fixed in the system. The whole of fixed nitrogen was utilized by plants grown. But, during composting, the fixed nitrogen was lost in the form of gases and leaching (reported by Dhar and Coworkers).

With forgoing discussions it can be concluded that direct application of organic matter both phosphated and unphosphated is superior to indirect application and at the same time higher doses of phosphorus was found to be superior to lower doses.

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# Study of Manganese Status of some deficient Medium Black Soils of Madhya Pradesh 

By

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#### Abstract

Surface soil samples $\left(0-9^{\prime \prime}\right)$ from 153 villages of four súb-divisions of Indore district viz., Mhow, Depalpur, Indore and Sawer which is a Manganese deficient area of the State, were analysed for exchangeable and reducible Manganese and for Cal. Carbonate, Organic Carbon and pH . The soils are found to be deficient in available Manganese but high in reducible Manganese. The soils of some of the villages in the locality were found to be acutely deficient in available Manganese.

Correlation studies indicated that available Manganese bears a negative relationship with pH and Cal. Carbonate but a positive relationship with Organic Carbon and reducible Manganese content of the soil. The correlation between available and reducible Manganese was found to be highly significant.

The importance of Micro-nutrients in Indian Agriculture can not be over emphasised especially in view of the efforts being made to improve Agriculture on scientific lines and on intensive cropping with high doses of fertilizers for increasing produgtion per unit area per unit time. In iecent years, the importance of the Micro nutrients in increasing Agricultural production and correcting deficiencies in plants : has been recognised by soil scientists all over the world. In the list of essential micronutrients, Manganese occupies a place second only to iron in order of importance for plant growth. Its availability directly depends on the soil $\mathrm{pH}, \mathrm{CaCO}_{3}$ and Organic Carbon content (Leeper, 1947. Biswas, 1951, Khanna, 1454). The medium black soils of India are rich in total Manganese but low in available Manganese (Bisuas 1951, 1953). It was also observed by earlier workers that the medium black cotton soil of Malwa plateau is deficient in available Manganese (Sharma and Motiramani 1964). The present study was conducted to find out the quantity of the two different forms of Manganese (Water Soluble + Exchangeable i. $t$. available and reducible Mn) in some medium black soils of Madhya Pradesh and to determine the relationship between different soil properties and reducible Manganese with available Manganese.


## Experimental

The surface soil samples were collected from the representative soils of 153 villages of Indore District. The samples were analysed for available Manganese (IN. Ammonium Acetate $\mathrm{pH} 7 \cdot 0$ ), reducible Manganese ( $\mathrm{IN} . \mathrm{NH}_{4}$ Acetate pH $7.0+0.2 \%$ hydroquinone), pH ( $1: 2.5$ soil water ratio by Beckman pH meter using glass electrode) as per procedures described by Jackson (1958) Organic Carbon (Walkley and Black's procedure, 1934) and Calcium Carbonate (rapid titration procedure as described by G. S. Piper, 1950). Estimation of Manganese was done by the Colorimetric procedure (Willard and Greathouse, 1917).

## Results and Discussion

The Manganese deficient medium black soils of M. P. which were taken up for study are clayey in texture with high calcium carbonate content. They are low in Organic Carbon content as shown in Table 1 (Muhr et al. 1965),

Table 1
Soil Characteristics (Average Values)

| S. N. | Sub-Division | No. of Samp. | pH | $\mathrm{CaCO}_{3} \%$ | O. C. \% |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 1. | Mhow | 60. | $8 \cdot 0$ | $25 \cdot 0$ | $0 \cdot 460$ |
| 2. | Depalpur | 55 | $8 \cdot 1$ | 23.0 | $0 \cdot 440$ |
| 3. | Indore | 53 | $8 \cdot 1$ | 23.0 | $0 \cdot 455$ |
| 4. | Sawer | 41 | $8 \cdot 2$ | $22 \cdot 0$ | $0 \cdot 445$ |
| Average of Total |  | 209 | $8 \cdot 1$ | $22 \cdot 0$ | 0.450 |

## Distribution of Manganese

Manganese status of soils is presented in Table 2. The available Manganese of these soils varied from 1.92 to 9.62 ppm with an average of 5.05 ppm for the district. Thus it can be safely said that these soils are poor in available Manganese (Toth-1951, Sharma and Motiramani, 1964).

Table 2
Distribution of Manganese (Available and Reducible)

| S. N. | Sub-Division | No. of Samps. | $\begin{array}{r} \text { Available } \mathrm{M} \\ \text { Range (ppm) } \end{array}$ | Manganese <br> Aver (ppm) | $\begin{gathered} \text { Reducible } \\ \text { Range (ppm) } \end{gathered}$ | Manganese <br> Aver (ppm) |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1. | Mhow | 60 | 1.92-9.62 | $6 \cdot 01$ | 21-16-481.0 | 255 |
| 2. | Depalpur | 55 | 2-16-5.20 | $5 \cdot 2$ | 48.10-481•0 | 258 |
| 3. | Indore | 53 | $2 \cdot 16-8 \cdot 42$ | $4 \cdot 7$ | 48.10-408.0 | 202 |
| 4. | Sawer | 41 | I•92-8.42 | $4 \cdot 4$ | 120.0-425.0 | 198 |

The easily reducible M. nganese content in the entire district ranges between $21 \cdot 16$ and 481 ppm with an average value of 236.43 ppm . Most of the samples showed a higher amount of easily reducible Manganese content (Leeper, 1947).

Table 3
Distribution of available Manganese Sub-division-wise

| Range of Mn (ppm) | Sub-Divisions |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | Mhow |  | alpur | Indore |  | Sawer |  |
|  | No. of Samp\% |  | No. of | Samp\% | No. | Samp\% | No. of | Samp\% |
| - 1.92-3.84 | 17. | 28\% | 7 | 12\% | 15 | 28\% | 16 | 39\% |
| $3 \cdot 84-6 \cdot 40$ | 17 | 28\% | 35 | 63\% | 33 | 64\% | 23 | 56\% |
| 6.40-9.6? | 26 | 44\% | 13 | 25\% | 5 | 8\% | 2 | 5\% |

Table 3 gives a comprehensive view of the distribution of samples under different ranges of percentages. Study of Table 3 reveals that all the samples have less than 10 ppm of available Manganese content showing thereby, that the entire area in Indore District is defictent in available Manganese content. It will be seen that more than $50 \%$ samples of each sub-divisiun have less than $6.4 \%$ Available Manganese. However, the soils of Sawer, Indore and Depalpur, subdivisions are comparatively more deficient than of Mhow sub-division. Sharma and Shinde (1968) h ive also reported similar results in black soils of Indore region. The results reported here also indicate that these soils are poor in available Manganese but high in reducible Manganese.

## Relationship of Available Manganese with reducible Manganese and other soil components

The relationship of Available Manganese with reducible Mangenese and other soils components such as $\mathrm{pH}, \mathrm{CaCO}_{3}$ content and Organic Carbon content is given below in Table 4 after working out the correlation coefficients.

Table 4
Correlation between Available Mangauese and different soil components

| S. No. | Correlation Between | r-Values |
| :--- | :--- | :---: |
| 1. | Available Manganese and soil pH | -0.12 |
| 2. | Available Manganese and $\mathrm{CaCO}_{3}$ content. | -0.129 |
| 3. | Available Manganese and Organic Carbon content. | +0.135 |
| 4. | Available Manganese and Reducible Manganese. | $+0.382^{*}$ |

*Significant at $1 \%$ level.
The correlation studies made between Available Manganese and other soil characteristics show that there is a negative correlation with pH and $\mathrm{CaCO}_{3}$ but a positive correlation with Organic Carbon and Reducible Manganese. However, the coefficient of correlation was found to be statistically significant, only in the case of Reducible Manganese.

The negative correlation in case of pH is in conformity with the observations of Berger (1948) Jones (1957), Zende et al. (1959) Randhawa et al. (1961), Mittal and Roy (1963) and also of previous workers in this laboratoiy viz; Sharma and Motiramani (1964) Dixit (1965), and Sharma (1966) who also found it to be highly significant statistically. This fact could be explained on the basis that at low pH , Manganese is present as di-valent Manganous ion which has higher solubility (Hewitt-1946), Berger and Gerloff-1948 and Biswas-1953) but wilh the increase in pH it gets converted into trivalent and tetravalent ions of Manganese which get precipitated as insoluble oxides and thus get converted into unexchangeable and therefore unavailable form (Sherman and Harmer-1942, Biswas-195"). Low values of exchangeable Manganese have also been reported in saline soils ol Rajasthan, Punjab and U. P. by Vinayak et al. (1964), Bhumbla and Dhingra (1964) and Agarwala et al. (1964) respectively.

The negative correlation observed in the case of $\mathrm{CaCO}_{3}$ percentage is also a harmony with the results reported by Yadav and Kalra (1964) Sherman and Fujimoto (1947) and Lohnis (1960) who reported decrease in available Manganese on addition of lime to Acid soils and also by previous workers in this
laboratory viz ; Sharma and Motiramani (1964), Dixit (1964) and Sharma (1966) who als? observed a negative significant correlation between available Manganese and $\mathrm{CaCO}_{3}$ content of soils. It is also in keeping with the negative relationship of Available Manganese with pH as pH is directly. related with $\mathrm{CaCO}_{3}$ content of soil.

The positive correlation observed in the case of Organic matter is not statistically significant. This is in conformity with the results of some earlier workers in this laboratory viz ; Sharma et al. (1964) and Dixit (1965) but is contradictory to the observations of Sharma (1966). In the literature also contradictory results have been reported by different workers (Kanwar and Randhawa 1967). Hence the role of Organic Matter in the conversion of Manganese oxides is not clear and there seems to be no significant relationship between Available Manganese and Organic Matter, especially when Organic Matter content is poor as is the case here. The positive correlation between Available and Reducible Manganese is also found to be highly significant statistically. This observation is also in agreement with the observations of earlier workers oiz., Sharma ot al. ' (1964), Dixit (1965) and Sharma (1966). Thus, there seems to exist a direct relationship between these two forms of Manganese, both of which are beneficial for plant growth (Jones and Leeper-1951).

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# Studies on Antibacterials Part IV Syntheses of Naphthalene analogues of Chalcone 

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Reactions of 4-hydroxy-3-nitro (or-5-bromo)-1-acetonaphthone with various aryl aldehydes in presence of $40 \%$ aqueous potassium hydroxide solution lead to the formation of coloured crystalline compounds. Some of these compounds have been found to be effective against $S$. aureus.

Chalcone compounds are known to possess antibacterial ${ }^{1,2}$, cardiovascular ${ }^{3}$, germicidal ${ }^{7}$, antifungal ${ }^{4}$, etc. activities. Davis and Eaton have suggessted ${ }^{5}$ that the compounds having ( $-\mathrm{CH}=\mathrm{CH}-\mathrm{CO}-$ ) grouping possess good insecticidal properties. In the present communication we have synthesized naphthalene analogues of chalcone in which phenyl and naphthyl groups are bridged by this grouping. These compounds have been prepared by the condensation of 4-hydroxy-3-nitro (or-3-bromo) -1- acetonaphthone with various aromatic aldehydes such as, bromo, chloro, methyl, hydroxy, nitro benzaldehydes and also with bromo veratricaldehydes, bromo vanillin, nitro vanillin, piperonal, bromo piperonal and 4-dimethylamino benzaldehyde in alcoholic medium ( $30-40 \mathrm{ml}$.), by using cold alkaline condensation.

It was observed that $50 \%$ alkali solution gave crystalline compounds with 4-hydroxy-3-nitro - 1 -acetonaphthone and the required aldehydes except nitro benzaldehydes which resulted into a mixture of unknown composition. Similar results were obtained in c.se of 3-bromo-4-hydroxy-1-n.phthyl substituted styryl ketones. In these cases the condensations were carried out with $10 \%$ potassium hydroxide solution at room temperature. This indicates that nitro benzaldehydes are sensitive towards the action of alkali in comparison with nitro methyl aryl ketones that is why with concentrated alkalis unidentifiable mixtures are obtained.

The chalcones of the 4-hydroxy-3-nitro (or 3-bromo) -1-acetonaphthone with, halo and nitro benzaldehydes were obtained in quite good yield (about $75-90 \%$ ), methyl and methoxy benzaldehydes gave about ( $45-60 \%$ ) of the theoretical yield, hydroxy benzaldehydes were obtained in very poor yield. The poor yield in case of hydroxy chalcones of naphthalene series may be on account of secondary reactions such as, resinification, etc. in presence of alkali solutions.

The presence of ketonic group in the synthesized compounds was characterized by preparing 2,4-dinitrophenylhydıazone derivatives. The analytical data and melting points are given in table 1.

The compounds of the 4-hydroxy-3-nitro-1-acetonaphthone series were screened for their antibacterial activity by observing their effect on $S$. aurcus. The activity was tested by using agar-cup method and the iesults were compared with benzoic acid (table \%).
4-hydroxy-3-nitro-1-naphthyl substituted styryl ketones


It was found that compounds having halogen atom or nitro group in the aldehyde component possessed marked antibacterial power. The balogen atom at ortho position w. r. s. to the conjugated system were more active than those at meta or para position.

## Experimental

All melting points are uncorrected.
3-bromo-4-hydroxy-1-acetonaphthone was prepared by bromination of 4-hydroxy-1-ace tonaphthone. ${ }^{6}$

## General Method for the Synthesis of chalcones

$40 \%$ aqueous solution of potassium hydroxide was gradually added with shaking to an ethanolic solution of equimolecular quantities. of 4-hydroxy-3-nitro (or 3-bromo)-1-acetonaphthone and the required aldehyde. The mixture was allowed to stand for a couple of hours with shaking at intervals. It was then poured into a beaker containing a mixture of crushed ice and water and then neutralized with $20 \%$ hydrochloric acid to a pH of 6 . The precipitated compound was collected on a Buchner funnel, washed with $1 \%$ sodium bicarbonate solution to eliminate the acid formed due to Cannizaro reaction and then finally from water. The crude compound on crystallization from a suitable solvent gave crystalline compounds (table 1).

TABLE 2
Antibacterial activity of 4-hydroxy-3-nitro-1-naphthyl substituted styryl ketones

| Substituents | Conc. in mg./ml. | Diam. of Zone of inhib. in mm | Substituents | Conc. in $\mathrm{mg} . / \mathrm{ml}$. | Diam. of Zone of inhib. in mm |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 2-bromo | 18 | 12 | 3-nitro | 20 | 10 |
| 4-bromo | 20 | 11 | 4-hydroxy- |  |  |
|  |  |  | 3-methoxy | 23 | $\bar{\square}$ |
| 3-chloro | 19 | 10 | 2-methoxy | 24 | 8 |
| 2, 6-dichloro | 18 | 12 | 3, 4-methylene- dioxy |  |  |
| 2-hydroxy | 24 | - | 4-dimethylamino | 22 | 6 |
| 4-hydroxy | 23 | - |  |  |  |
| 2, 4-dihydroxy | 25 | - |  |  |  |
| 2-methoxy | 27 | - |  |  |  |
| 2, 3-dimethoxy | 26 | - |  |  |  |
| Benzoic acid | 20 | 12 |  |  |  |

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# Intrinsic properties of steady magnetogasdynamic 

## flows II

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#### Abstract

Herein we have studied the variations of flow parameters for plane hydromagnetic flows along the streamlines and their orthogonal trajectories; considering two sy stems of plane curves one related to the streamlines and the other to their orthogonal trajectories.


## 1. Introduction

Warsic [5] in his recent paper has studied the variation of flow parameters for steady plane magnetogasdynamic flows. Defining two systems of plane orthogonal curves of congruences one related to the magnetic lines of force in earlier paper [4] we have studied the intrinsic properties of plane hydromagnetic flows. In this paper we have studied the variations of flow parameters along the streamlines and their orthogonal trajectories for plane hydromagnetic flows. Decomposing the equation governing conservation of mass into intrinsic form, it is proved that the streamline patterns as a system of concentric circles or parallel lines, if the momentum per unit mass of fluid is conserved along a streamline, which is independent of the magnetic field. It is also observed that the coincidence of isovels with the streamlines does not imply the coincidence of isobars with them as in the case of non-magnetic flows. The field equation is independent of the tangentral component of the magnetic field to the streamline. Introducing the velocity of the sound, we have transformed the equation into intrinsic form and observed that the variation of mach number along a streamline is equal to the product of the curvature of normal congruence and the mach number for Chaplygon's gas.
2. (A) Basic Equalions

The basic equations governing steady hydromagnetic flows, in the absence of extraneous forces are given below in the usual notation [2]:

$$
\begin{equation*}
\operatorname{div}(\overrightarrow{\rho q})=0 \tag{1}
\end{equation*}
$$

$$
\begin{equation*}
\rho(\vec{q} \cdot \nabla) \vec{q}=-\nabla p+\mu e \vec{J} \lambda \overrightarrow{\vec{H}} \tag{2}
\end{equation*}
$$

$$
\begin{equation*}
(\vec{q} \cdot \nabla) \vec{H}+\vec{H} \operatorname{div} \vec{q}-(\vec{H} \cdot \nabla) \vec{q}=0 \tag{3}
\end{equation*}
$$

$$
\begin{equation*}
p=\rho^{\gamma}{ }_{e} S_{i}^{H} \lambda_{c v} \tag{4}
\end{equation*}
$$

$$
\begin{equation*}
\operatorname{div} \vec{H}=0 \tag{5}
\end{equation*}
$$

$$
\begin{equation*}
\overrightarrow{(q} \cdot \nabla) S=0 \tag{6}
\end{equation*}
$$

where $\vec{q}, \vec{H}, \rho, p, S, \mu_{e}, \gamma, \lambda$ and $\sigma_{v}$ are the velocity vector, the magnetic field, the density, the pressure, the specific entropy. the magnetic permeability, the adiabatic exponent, Joul's'constant and the specific heat at constant volume respectively.
(B) Geometric Results

Considering $\vec{t}$ and $\vec{n}$ as two plane curves of congruence, one of them related to the streamlines and the other to their orthogonal trajectories and also denoting $\left(\frac{d}{d s}, \frac{d}{d n}\right)$ as the directional derivatives along these vectors, we have the fallowing geometric results [3]:

$$
\begin{equation*}
\frac{\overrightarrow{d \gamma}}{d s}=\vec{t}=\frac{\vec{q}}{q} \tag{7}
\end{equation*}
$$

$$
\begin{array}{ll}
\frac{\vec{d} t}{d s}=\vec{n} k & \text { (9) } \frac{d \vec{n}}{d s}=-\vec{t} k \\
\frac{\overrightarrow{d n}}{d n}=+\vec{t} k^{\prime} & \text { (11) } \frac{d \vec{t}}{d n}=-\vec{n} k^{\prime}
\end{array}
$$

where $k$ and $k^{\prime}$ are the curvatures of the streamlines and their orthogonal trajectories respectively.

Also using the irrotational property of ' $\nabla f^{\prime}$ ' we obtain

$$
\begin{equation*}
k^{\prime} \frac{d f}{d n}+\frac{d}{d n}\left(\frac{d f}{d s}\right)=k \frac{d f}{d s}+\frac{d}{d s}\left(\frac{d f}{d n}\right) \tag{12}
\end{equation*}
$$

This has to be satisfied by any flow quantity.
3. Decomposition into Intrinsic Form:

In this section, we shall make use of the above geometric results to transform the basic equations into intrinsic form and study some of the kinetic and kinematic properties of fluid flows.

Using (7) in (1) we have

$$
\begin{equation*}
k^{\prime}=\frac{d}{d s} \log (\rho q) \tag{13}
\end{equation*}
$$

which gives the curvature of orthogonal trajectories of the streamlines. If the momentum of a unit fluid mass is conserved along an individual streamline, then the stream patterns must be either a system of concentric circles or a family of parallel straight lines, i.e., the flow due to a vortex type or rectilinear.

Forming the scalar product of (2) by $\vec{t}$ and $\vec{n}$, we obtain successively

$$
\begin{align*}
\rho q \frac{d q}{d s} & =-\frac{d p}{d s}-\mu_{e} J H_{n}  \tag{14}\\
\rho q^{2} k & =-\frac{d p}{d n}+\mu_{e} J H_{t} \tag{I5}
\end{align*}
$$

where $H_{t}$ and $H_{n}$ are the resolved parts of the magnetic field along the streamlines and their orthogonal trajectories. From (14) we observe that the coincidence of either isovels or isobars with the streamlines does not imply the remaining as in the case of nonmagnetic flows. c.f. Nemenyi and Prim [1]. Also it is evident from (14) that the coincidence of any two of isovels, isobars and the magnetic lines with the streamlines implies the remaining. From (15) we obtain the current density expression, when the streamlines and isobars intersect orthogonally, i.e., the pressure remains uniform along the orthogonal trajectories of the streamlines as

$$
\begin{equation*}
J=\frac{\rho k q^{2}}{\mu_{e} H_{l}^{-}} \tag{16}
\end{equation*}
$$

Also if the streamlines are straight, the magnetic lines are either circles or straight lines.

The field equation (3) in intrinsic form can be decomposed as

$$
\begin{gather*}
q\left(\frac{d H_{t}}{d s}-k H_{n}\right)+H_{t}\left(\frac{d q}{d s}-k q\right)  \tag{17}\\
-\left(H_{t} \frac{d q}{d s}+H_{n}\left(\frac{d q}{c^{\prime} n}\right)=0\right. \\
q\left(k H_{t}+\frac{d H_{n}}{d s}\right)+H_{n}\left(\frac{d q}{d s}-k^{\prime} q\right)  \tag{18}\\
-\left(k q H_{t}-k^{\prime} q H_{n}\right)=0
\end{gather*}
$$

The conservation of magnetic field can be written as

$$
\begin{equation*}
\frac{d H_{t}}{d s}+\frac{d H_{n}}{d n}-\left(k H_{n}+k^{\prime} H_{t}\right)=0 \tag{19}
\end{equation*}
$$

Eliminating $\frac{d H_{t}}{d s}$ from (17) and (19) we obtain

$$
\begin{equation*}
\frac{d}{d s}\left(q H_{n}\right)=0 \tag{20}
\end{equation*}
$$

Also eliminating $\frac{d H_{n}}{d s}$ from (18) and (19) we obtain

$$
\begin{equation*}
\frac{d}{d n}\left(q H_{n}\right)=0 \tag{21}
\end{equation*}
$$

Eliminating ' $q$ ' and $H_{n}$ from (20) and (21) we obtain

$$
\frac{d H_{n}}{d s} \frac{d q}{d n}-\frac{d q}{d s} \frac{d H_{n}}{d n}=0
$$

This is the field equation in intrinsic form
Using (4), (2) can be written as

$$
\begin{equation*}
(\vec{q} \cdot \nabla) \vec{q}+\frac{2 c}{\gamma-1} \nabla c=\frac{c^{2}}{\lambda c_{p}(\gamma-1)} \nabla S+\frac{\mu_{e}}{\rho} \vec{J} \lambda \vec{H} \tag{23}
\end{equation*}
$$

Forming scalar product of (23) by $\vec{i}, \vec{n}$ and using (4) we get

$$
\begin{align*}
& q \frac{d q}{d s}+\frac{2 c}{\gamma-1} \frac{d c}{d s}=-\frac{\mu_{e}}{\rho} J H_{n}  \tag{24}\\
& k q^{2}+\frac{2 c}{\gamma-1} \frac{d c}{d n}=\frac{c^{2}}{\lambda c_{p}(\gamma-1)} \frac{d S}{d n}+\frac{\mu_{e}}{\rho} J H_{t} \tag{25}
\end{align*}
$$

From (14) and (24) we conclude that the coincidence of any two of the isovels, isobars, magnetic lines and sound vels with the streamlines implies the remaining two.

Also using (4) in (1) we obtain

$$
\begin{equation*}
k^{\prime}=\frac{d}{d s} \log M+\frac{\gamma+1}{\gamma-1} \frac{d}{d s} \log c \tag{26}
\end{equation*}
$$

Eliminating the density from (1) and (4) we obtain

$$
\begin{equation*}
k^{\prime}=\frac{d}{d s} \log \left(\frac{2}{c^{\gamma-1}} q\right) \tag{27}
\end{equation*}
$$

For Chaplygin's gas (27) simplifies to

$$
\begin{equation*}
k^{\prime}=\frac{d}{d s} \log M \tag{28}
\end{equation*}
$$

This gives the variation of the mach number along a streamline.

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# Kinetics of oxidation of Toluene by Peroxydisulphate Ion 

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#### Abstract

The oxidation of toluene has been studied in sulphuric acid medium. The reaction is first order w. r. t. peroxydisulphate and zero order w. r. t. toluene. The specific rate is related to $\mathrm{AgNO}_{3}$ concentration by the relationship : $$
k=0.06 \times 10^{-3}+1.3 \mathrm{C}_{\mathrm{Ag}^{+}}{ }^{+}
$$

The salt effect is negative. The energy parameters are $\Delta \mathrm{E}=12 \cdot 4 \mathrm{Kcal} /$ mole, $A=7.67 \times 10^{3}$ litre mole ${ }^{-1} \mathrm{sec}^{-1}$ and $\Delta S=-42.85 \mathrm{E}$. U.


The uncatalysed oxidation of different organic substrates by peroxydisulphate ion is generally slow. They are appreciably catalysed by silver ion. The kinetic studies involving peroxydisulphate ion as an oxidising agent have been recently reviewed by House ${ }^{1}$. Bacon and Doggai $t^{2}$ found that toluene on oxidation by peroxydisulphate ion yields a number of products-benzoic acid, benzaldehyde, bibenzyl and a small amount of resinous products. However, toluene on oxidation with permangnate in aqueous medium yields benzoic acid as the chief product ${ }^{3}$, but $\mathrm{V}^{\mathrm{v}}$ oxidation of toluene gives benzaldehyde as the main product ${ }^{4}$. A review of the literature, however, shows that the kinetic study of the oxidation of toluene by peroxydisulphate ion has not been investigated so far. The present paper deals with the kinetic study of the $\mathrm{Ag}^{+}$catalysed oxidation of toluene by peroxydisulphate ion.

## Experimental

The experimental procedure is similar to that adopted in the studies on phthalic acid-peroxydisulphate reaction ${ }^{5}$. All chemicals used were of Analar grade and toluene of Analar grade was used after redistillation.

The reaction was slow but the rate was measurable after addition of $\mathrm{AgNO}_{3}$ at $45^{\circ} \mathrm{C}$ and above.

## Results and Measurements

Since the self decomposition of peroxydisulphate ion at this temperature is appreciable in the presence of $\mathrm{AgNO}_{3}$, the rate of self decomposition was studied simultaneously and the rate of disappearance of $\mathrm{S}_{2} \mathrm{O}_{8}{ }_{8}$ due to toluene was evaluated by subtracting the rate constant of self decomposition from the observed rate. To avoid precipitate and emulsion formation, the reaction was carried out in sulphuric acid medium ( 0.025 M ).

$$
\begin{aligned}
& \text { FigI } \\
& \begin{aligned}
\text { Ag.NO3 } & =0.001 \mathrm{M} \\
\text { HzSO4 } & =0.025 \mathrm{M} \\
\text { Toluene } & =0.057 \mathrm{M}
\end{aligned}
\end{aligned}
$$



The effect of initial concentration of $\mathrm{K}_{2} \mathrm{~S}_{2} \mathrm{O}_{8}$ on its rate of disappearance was determined at constant concentrations of toluene, $\mathrm{AgNO}_{3}$ and $\mathrm{H}_{2} \mathrm{SO}_{4}$, vide Table 1.

Table 1

$$
\begin{array}{ll}
\mathrm{AgNO}_{3}=0.001 \mathrm{M} & \text { Toluene }=0.057 \mathrm{M} \\
\mathrm{H}_{2} \mathrm{SO}_{4}=0.025 \mathrm{M} & \text { Temp. }=45^{\circ} \mathrm{C} . \\
\hline
\end{array}
$$

| S. No. | Concn. of <br> $\mathrm{K}_{2} \mathrm{~S}_{2} \mathrm{O}_{8}(\mathrm{M})$ | $\mathrm{k}^{\prime} \times 10^{3}$ <br> $\mathrm{~min}^{-1}$ | $\mathrm{k}^{n} \times 10^{3}$ <br> $\mathrm{~min}^{-1}$ | $\left(\mathrm{k}^{n}-\mathrm{k}^{\prime}\right) \times 10^{3}$ <br> $\mathrm{~min}^{-1}$ |
| :---: | :---: | :---: | :---: | :---: |
| 1 | 0.005 | 1.76 | 3.90 | 2.14 |
| 2 | 0.010 | 1.92 | 3.44 | 1.52 |
| 3 | 0.015 | 2.01 | 3.34 | 1.33 |
| 4 | 0.020 | 2.19 | 3.11 | 0.92 |
| 5 | 0.025 | 2.23 | 2.79 | 0.56 |
| 6 | 0.030 | 2.24 | 2.73 | 0.49 |

$k^{\prime \prime}=$ Specific rate in presence of toluene
$\mathrm{k}^{\prime}=$ Specific rate in absence of toluene
It is seen that at all concentrations of peroxydisulphate ion studied, the reaction is first order. However, the rate constant is found to decrease with the increase in $\mathrm{S}_{2} \mathrm{O}_{8}{ }^{n}$ concentration. This decrease could be ascribed due to increase in ionic strength but this is noi possible because in any particular run the rate constant does not decrease with time although the ionic strength o the reaction mixture increases as the reaction proceeds. The other possibility for this behaviour is that the $\mathrm{K}^{+}$exerts a specific inhibitory effect. This type of behaviour has also been shown to be present in other peroxydisulphate redux reactions, $C . f$. Khulbe and Srivastava's $s^{8}$ observations on the oxidation of glucose by $\mathrm{K}_{2} \mathrm{~S}_{2} \mathrm{O}_{8}$.

## Effect of toluene concentration:

The rate constants corresponding to different initial concentration of toluene are given in Table 2.

Table 2

$$
\begin{array}{ll}
\mathrm{K}_{2} \mathrm{~S}_{2} \mathrm{O}_{8}=0.010 \mathrm{M} & \mathrm{AgNO}_{3}=0.001 \mathrm{M} \\
\mathrm{H}_{2} \mathrm{SO}_{4}=0.025 \mathrm{M} & \text { Temp. }=45^{\circ} \mathrm{C} \\
\hline
\end{array}
$$

| Concn. of <br> toluene <br> $(\mathrm{M})$ | 0.038 | 0.057 | 0.076 | 0.114 | 0.142 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\left(\mathrm{k}^{\prime \prime}-\mathrm{k}^{\prime}\right) \times 10^{3}$ <br> $\mathrm{~min}^{-1}$ | 1.51 | 1.52 | 1.68 | 1.55 | 1.67 |

$$
\mathbf{k}^{\prime}=1.92 \times 10^{-3} \mathrm{~min}^{-1}
$$



It is clear from the above table that the rate constant is not appreciably affected by a change in toluene concentration.

Since the reaction follows a first order behaviour when the concentration of toluene is in excess (vide kinetic runs 1 and 4 and 5 under Tables 1 and 2 respectively) and also when the concentration of $\mathrm{S}_{2} \mathrm{O}_{8}{ }^{2-}$ and toluene are similar, the reaction, therefore, is first order in $\mathrm{S}_{2} \mathrm{O}_{8}{ }^{2-}$ and zero order in toluene.
Effect of $\mathrm{AgNO}_{3}$ concentration:
Table 3 summarises the results for the effect of $\mathrm{AgNO}_{3}$ concentration on the rate of oxidation of toluene.

Table 3

|  | $\begin{aligned} & \mathrm{K}_{2} \mathrm{~S}_{2} \mathrm{O}_{8}=0.010 \mathrm{M} \\ & \text { Toluene }=0.057 \mathrm{M} \end{aligned}$ |  | $\begin{aligned} & \mathrm{H}_{2} \mathrm{SO}_{4}=0.025 \mathrm{M} \\ & \mathrm{Temp} . \end{aligned}=45^{\circ} \mathrm{C} \text {. }$ |  |
| :---: | :---: | :---: | :---: | :---: |
| S. No. | $\mathrm{AgNO}_{3}$ conen. <br> (M) | $\begin{aligned} & \mathrm{k}^{\prime} \times 10^{3} \\ & \mathrm{~min}^{-1} \end{aligned}$ | $\begin{gathered} \mathrm{k}^{7} \times 10^{3} \\ \min ^{-1} \end{gathered}$ | $\begin{gathered} \left(\mathrm{k}^{\prime \prime}-\mathrm{k}^{\prime}\right) \times 10^{3} \\ \min ^{-1} \end{gathered}$ |
| 1 | 0.0005 | $1 \cdot 38$ | $2 \cdot 06$ | 0.68 |
| 2 | 0.0010 | $1 \cdot 92$ | $3 \cdot 44$ | 1.52 |
| 3 | 0.0015 | $3 \cdot 19$ | 4.98 | 1.79 |
| 4 | $0 \cdot 0020$ | $4 \cdot 00$ | 6.08 | $2 \cdot 08$ |
| 5 | $0 \cdot 0025$ | $4 \cdot 83$ | $8 \cdot 28$ | $3 \cdot 45$ |

The specific rate is found to be linearly related to the concentration of $\mathrm{AgNO}_{3}$ and is given by the expression :

$$
k=0.06 \times 10^{-3}+1.3 \mathrm{C}_{\mathrm{Ag}}+
$$

Effect of temperature :
Results for kinetic runs at five different tempertures are tabulated in Table 4.
Table 4

$$
\begin{array}{ll}
\mathrm{K}_{2} \mathrm{~S}_{2} \mathrm{O}_{8}=0.010 \mathrm{M} & \mathrm{AgNO}_{3}=0.001 \mathrm{M} \\
\mathrm{H}_{2} \mathrm{SO}_{4}=0.025 \mathrm{M} & \text { Toluene }=0.057 \mathrm{M}
\end{array}
$$

| S. No. | Temp. ${ }^{\circ} \mathrm{A}$ | $\underset{\mathbf{m i n}^{-1}}{\mathbf{k}^{\prime} \times 10^{3}}$ | $\begin{array}{r} \mathrm{k}^{\prime \prime} \times 10^{3} \\ \min ^{-1} \end{array}$ | $\underset{\min ^{-1}}{\left(k^{\prime \prime}-k^{\prime}\right) \times 10^{3}}$ | Temp. coefficient | $\triangle \mathrm{E}$ <br> Kcals/ mole | $\begin{array}{r} \mathrm{A} \times 10 \\ \text { litre } \\ \text { mole } \\ \mathrm{sec}^{-} \end{array}$ | $\triangle \mathrm{S}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 308 | $1 \cdot 12$ | $1 \cdot 91$ | 0.79 |  |  | $7 \cdot 48$ | $-42 \cdot 87$ |
| 2 | 313 | $1 \cdot 48$ | $2 \cdot 64$ | -16 | $1 \cdot 92$ | $12 \cdot 7$ | $7 \cdot 89$ | -42.76 |
| 3 | 318 | $1 \cdot 92$ | $3 \cdot 44$ | 52 | 1.85 | $12 \cdot 2$ | 7.61 | -42.84 |
| 4 | 323 | $2 \cdot 60$ | $4 \cdot 72$ | $2 \cdot 12$ | $1 \cdot 80$ | $12 \cdot 3$ | $7 \cdot 82$ | -42:86 |
| 5 | 328 | $3 \cdot 69$ | 6. 44 | 2.75 |  |  | 7.56 | -42.95 |

$\Delta \cdot E$ for self decomposition $=12 \cdot 1 \mathrm{~K}$ cals mole ${ }^{-1}$
[ 479 ]


GRAPH 2

A perusal of the above values shows that the energy of activation for themal decomposition of $\mathrm{S}_{2} \mathrm{O}_{8}{ }^{\prime \prime}$ and for toluene oxidation is of the same order suggesting that the mechanism of both the processes is similar.

A large negative value of $\Delta S$ suggests that the activated complex formed with the reductant is rather rigid in structure and further the rate determining reaction is between two oppositely charged ions. It may be mentioned here that a negative value for the entropy of activation has also been reported for other $\mathrm{Ag}^{+}$ catalysed redox reactions of peroxydisulphate involving organic substrates $C . f$. Khulbe and Srivastava ${ }^{7}$.

## Salt effect:

In order to determine the effect of change of ionic strength on the rate, the reaction was carried out in the presence of different concentrations of $\mathrm{K}_{2} \mathrm{SO}_{4}$, added as neutral salt.

Table 5

$$
\begin{array}{rll}
\mathrm{K}_{2} \mathrm{~S}_{2} \mathrm{O}_{8}=0.010 \mathrm{M}, & \mathrm{AgNO}_{3}=0.001 \mathrm{M}, & \text { Temp. }=45^{\circ} \mathrm{C} \\
\mathrm{H}_{2} \mathrm{SO}_{4}=0.025 \mathrm{M}, & \text { Toluene }=0.057 \mathrm{M} &
\end{array}
$$

| S. No. | $\mathrm{K}_{2} \mathrm{SO}_{4}$ concn. (M) | $\begin{gathered} \mathbf{k}^{\prime} \times 10^{2} \\ \min ^{-1} \end{gathered}$ | $\underset{\min ^{-1}}{\mathrm{k}^{n} \times 10^{3}}$ | $\underset{\min ^{-1}}{\left(\mathrm{k}^{\prime}-\mathrm{k}^{\prime}\right) \times 10^{3}}$ | $\log \left(k^{\prime \prime}-k^{\prime}\right)$ | $\mu$ | $\mu^{\frac{2}{2}}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | 0.000 | 1.92 | $3 \cdot 44$ | $1 \cdot 52$ | -2.817 | $0 \cdot 106$ | 0.326 |
| 2 | 0.010 | $1 \cdot 80$ | $3 \cdot 02$ | $1 \cdot 22$ | -2.914 | $0 \cdot 136$ | 0.369 |
| 3 | 0.020 | $1 \cdot 75$ | $2 \cdot 81$ | $1 \cdot 06$ | -2.975 | $0 \cdot 166$ | $0 \cdot 407$ |
| 4 | 0.030 | $1 \cdot 70$ | $2 \cdot 71$ | $1 \cdot 01$ | -2.994 | $0 \cdot 196$ | 0.443 |
| 5 | 0.040 | 1.69 | $2 \cdot 65$ | 0.96 | -3.019 | $0 \cdot 226$ | $0 \cdot 475$ |
| 6 | 0.050 | $1 \cdot 67$ | $2 \cdot 51$ | 0.84 | -3.079 | 0.256 | 0.502 |
| 7 | 0.075 | $1 \cdot 41$ | $2 \cdot 21$ | 0.80 | -3.097 | 0.331 | 0.575 |
| 8 | $0 \cdot 100$ | $1 \cdot 13$ | 1-89 | 0.76 | -3.114 | 0.406 | $0 \cdot 637$ |

It is seen that the salt effect is negative. However, neither ( $k^{\prime \prime}-k^{\prime}$ ) nor log ( $k^{\prime \prime}-k^{\prime}$ ) decrease linearly with the increase of $\mu$ or $\mu^{\frac{1}{2}}$, therefore. it is not possible to decide the exact nature of the salt effect. It is worth mentioning here that since the reaction has to be carried out in sulphuric acid medium $(0.025 \mathrm{M})$, Brönsted relationship for salt effect is naturally not applicable.

Thus the reaction between potassium peroxydisulphate and toluene is first order w.r.t. $\mathrm{S}_{2} \mathrm{O}_{8}{ }^{n}$, zero order w.r.t. toluene and negative salt effect for the $\mathrm{Ag}{ }^{+}$ catalysed reaction suggests that the rate determining step may be between two oppositely charged ions.

Further work to study the kinetics of different stages and to isolate and identify different intermediate products is in progress after that only, a detailed reaction mechanism can be put forward.

## Acknowledgements

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## CONTENTS

Page
Strong almost convergence of order ( $\mathbf{r}$ ) of the Fourier series

- Sushil Kumar Srivastava ..... 355
Some Theorems on Transform Umesh Chandra Sharma ..... 359
Convolution in a General Transform D. G. Joshi367
Superposition of Perfect Fluid and Electromagnetic Fields in Spheri. cally Symmetric Space.Time in General Relativity . K. P. Singh and M. C. Srivastava ..... 372
On the Rate of Decay of Solutions of Neutral Functional Differential
On the Rate of Decay of Solutions of Neutral Functional Differential Systems B. G. Pachpatte ..... 383
On order of partial sum of Laplace series S. K. Tiwari ..... 389
Effect of Viscosity on the Stability of Superposed Fluids S. N. Verma and H. G. Khare ..... 295
Series Involving Legendre Polynomials R. C. Verma ..... 399
A note on certain bilinear and bilateral generating relations for the Laguerre Polynomial H. M. Srivastava
407
407
Reactions of Ortho-esters of Germanium-reactions of ethyl Orthoger. manate with acetic and Benzoic acids R. C. Mehrotra and G. Chandra ..... 415
On Resultant Kernel Functions . O. P. Sharma ..... 421
Application of generalized function method to the theory water waves Lokenath Debnath ..... 431
Kinetics and Mechanism of Reaction between Bromine and Quinoline in Carbon-Tetrachloride Solution Part II . . V. N. P. Srivastava ..... 440
Effect of algae and Phosphorus on the formation of aminoacids
Oxidation of a-Hydroxy acids by Vanadium (V) - R. K. Prasad andEffect of Direct and Indirect application of Organic Matter (Phosphat-ed and unphosphated) on the yield of wheat cropN. R. Dhar, S. S. Singh, R. P. S. Chauhan and H. G. Sharma459
Study of Manganese Status of some deficient Medium Black Soils of Madhya Pradesh . . . . . S. C. Sharma and S. B. Sinha ..... 463
Studies on Antibacterial Part IV Syntheses of Naphthalene analogues of Chalcone R. S. Tewari and S. S. Misra ..... 468
Intrinsic properties of steady magnetogasdynamic flows II ..... 471
Kinetics of oxidation of Toluene by Peroxydisulphate Ion
S. P. Srivastava and S. K. Singhal ..... 475


[^0]:    *Presented to the thirty-eighth Annual Session of the National Academy of Sciences of India.
    $\dagger$ It may be of interest to note that the formulas (3). and (4) appear incorrectly on page 938 of Jain's paper referred to above.

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